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## 0 Measure Theory

### 0.1 Key Theorems

Theorem 0.1. Fubini's Theorem
http://en.wikipedia.org/wiki/Fubini\'s_theorem

Suppose $A$ and $B$ are complete measure spaces. Suppose $f(x, y)$ is $A \times B$ measurable. If

$$
\int_{A \times B}|f(x, y)| d(x, y)<\infty
$$

where the integral is taken with respect to a product measure on the space over $A \times B$, then

$$
\int_{A}\left(\int_{B} f(x, y) d y\right) d x=\int_{B}\left(\int_{A} f(x, y) d x\right) d y=\int_{A \times B} f(x, y) d(x, y)
$$

the first two integrals being iterated integrals with respect to two measures, respectively, and the third being an integral with respect to a product of these two measures.

## Corollary:

If $f(x, y)=g(x) h(y)$ for some functions $g$ and $h$, then

$$
\int_{A} g(x) d x \int_{B} h(y) d y=\int_{A \times B} f(x, y) d(x, y)
$$

the third integral being with respect to a product measure.

Theorem 0.2. Tonelli's Theorem
http://en.wikipedia.org/wiki/Fubini\'s_theorem\#Tonelli.27s_theorem

Suppose that $A$ and $B$ are $\sigma$-finite measure spaces, not necessarily complete. If either

$$
\int_{A}\left(\int_{B}|f(x, y)| d y\right) d x<\infty \text { or } \int_{B}\left(\int_{A}|f(x, y)| d x\right) d y<\infty
$$

then

$$
\int_{A \times B}|f(x, y)| d(x, y)<\infty
$$

and

$$
\int_{A}\left(\int_{B} f(x, y) d y\right) d x=\int_{B}\left(\int_{A} f(x, y) d x\right) d y=\int_{A \times B} f(x, y) d(x, y)
$$

Tonelli's theorem is a successor of Fubini's theorem. The conclusion of Tonelli's theorem is identical to that of Fubini's theorem, but the assumptions are different. Tonelli's theorem states that on the product of two -finite measure spaces, a product measure integral can be evaluated by way of an iterated integral for nonnegative measurable functions, regardless of whether they have finite integral. A formal statement of Tonelli's theorem is identical to that of Fubini's theorem, except that the requirements are now that $(X, A, \mu)$ and $(Y, B, \nu)$ are $\sigma$-finite measure spaces, while $f$ maps $X \times Y$ to $[0, \infty]$.

Theorem 0.4. Hölder's Inequality
Theorem 12.54 on page 356

Let $1 \leq p, q \leq \infty$ and $\frac{1}{p}+\frac{1}{q}=1$. If $f \in L^{p}(X, \mu)$ and $g \in L^{q}(X, \mu)$, then $f g \in L^{1}(X, \mu)$ and

$$
\left|\int f g d \mu\right| \leq\|f\|_{p}\|g\|_{q}
$$

Theorem 0.5. Young's Inequality
Theorem 12.58 on page 359
Let $1 \leq p, q, r \leq \infty$ and $\frac{1}{p}+\frac{1}{q}=1+\frac{1}{r}$. If $f \in L^{p}\left(\mathbb{R}^{n}\right)$ and $g \in L^{q}\left(\mathbb{R}^{n}\right)$, then $f * g \in L^{r}\left(\mathbb{R}^{n}\right)$ and

$$
\|f * g\|_{r} \leq\|f\|_{p}\|g\|_{q}
$$

Theorem 0.6. Lebesgue Dominated Convergence Theorem
Theorem 12.35 on page 348

Suppose that $\left(f_{n}\right)$ is a sequence of integrable functions, $f_{n}: X \rightarrow \overline{\mathbb{R}}$, on a measure space $(X, \mathcal{A}, \mu)$ that converges pointwise to a limiting function $f: X \rightarrow \overline{\mathbb{R}}$. If there is an integrable function $g: X \rightarrow[0, \infty]$ such that

$$
\left|f_{n}(x)\right| \leq g(x) \quad \forall x \in X, n \in \mathbb{N}
$$

then $f$ is integrable and

$$
\lim _{n \rightarrow \infty} \int f_{n} d \mu=\int f d \mu
$$

Theorem 0.7. Cauchy-Schwarz Inequality http://en.wikipedia.org/wiki/Cauchy-Schwarz_inequality

Formal Statement: For all vectors $x, y$ of an inner product space,

$$
\begin{aligned}
|\langle x, y\rangle|^{2} & \leq\langle x, x\rangle\langle y, y\rangle \\
|\langle x, y\rangle| & \leq\|x\|\|y\|
\end{aligned}
$$

Square of a Sum:

$$
\left|\sum_{i=1}^{n} x_{i} y_{i}\right|^{2} \leq \sum_{i=1}^{n}\left|x_{i}\right|^{2} \sum_{i=1}^{n}\left|y_{i}\right|^{2}
$$

In $L^{2}$ :

$$
\left|\int f(x) g(x) d x\right|^{2} \leq \int|f(x)|^{2} d x \int|g(x)|^{2} d x
$$

## 7 Fourier Series

### 7.1 Fourier Series

## Definition 7.1. $2 \pi$-periodic

page 149

A function $f: \mathbb{R} \rightarrow \mathbb{C}$ is $2 \pi$-periodic if

$$
f(x+2 \pi)=f(x) \quad \forall x \in \mathbb{R}
$$

A $2 \pi$-periodic function may be indentified with a function on the unit circle, or one-dimensional torus, $\mathbb{T}=\mathbb{R} /(2 \pi \mathbb{Z})$. The space $C(\mathbb{T})$ is the space of continuous functions from $\mathbb{T}$ to $\mathbb{C}$, and $L^{2}(\mathbb{T})$ is the completion of $C(\mathbb{T})$ with respect to the $L^{2}$-norm,

$$
\|f\|=\left(\int_{\mathbb{T}}|f(x)|^{2} d x\right)^{1 / 2}
$$

$L^{2}(\mathbb{T})$ is a Hilbert space with respect to the inner product

$$
\langle f, g\rangle=\int_{\mathbb{T}} \overline{f(x)} g(x) d x
$$

## Definition 7.2. $L^{p}(\mathbb{T})$

page 92 and Notes $1 / 3 / 11$
$L^{p}(\mathbb{T}):=$ the space of Lebesgue measurable functions, $f: \mathbb{T} \rightarrow \mathbb{C}$ such that $\int_{\mathbb{T}}|f|^{p} d x<\infty$, where $1 \leq p<\infty$. We define the $L^{p}$-norm as:

$$
\|f\|_{p}=\left(\int_{\mathbb{T}}|f|^{p} d x\right)^{1 / p}
$$

For $p=\infty, L^{\infty}(\mathbb{T})$ is the space of Lebesgue measurable functions that are essentially bounded on $\mathbb{T}$, meaning that $f$ is bounded on every subset of $\mathbb{T}$ with nonzero measure. The norm on $L^{\infty}(\mathbb{T})$ is the essential supremum

$$
\|f\|_{\infty}=\inf \{M| | f(x) \mid \leq M \text { a.e. in } \mathbb{T}\}
$$

We identify $f$ with $g$ if $f=g$ a.e. (almost everywhere, except possibly on a set of measure 0 ).

## Theorem 7.3.

Notes 1/3/11
$L^{p}(\mathbb{T})$ with the norm $\|f\|_{L^{p}}=\left(\int_{\mathbb{T}}|f|^{p} d x\right)^{1 / p}$ is a Banach space.

## Theorem 7.4.

Notes $1 / 3 / 11$
$C(\mathbb{T})$ is dense in $L^{p}(\mathbb{T})$ for $1 \leq p<\infty$.
Note: $C(\mathbb{T}):=$ the space of continuous functions $f: \mathbb{T} \rightarrow \mathbb{C}$

## Proposition 7.5.

Notes $1 / 3 / 11$
$p>q \Rightarrow L^{p}(\mathbb{T}) \subset L^{q}(\mathbb{T}) \quad$ and $\quad\|\cdot\|_{p} \geq\|\cdot\|_{q}$
Also,
$L^{1}(\mathbb{T}) \supset L^{2}(\mathbb{T}) \supset \ldots \supset C(\mathbb{T})$

## Example 7.6. Fourier Basis Example

Notes 1/3/11

$$
\begin{aligned}
& \sum_{n \neq 0} \frac{1}{|n|} e^{i n x}=f(x) \\
& \sum_{n \neq 0} \frac{1}{|n|^{2}}=2 \sum_{n=1}^{\infty} \frac{1}{n^{2}}<\infty \\
& \lim _{N \rightarrow \infty} \int\left|f(x)-\sum_{n=-N, n \neq 0}^{N} \frac{1}{|n|} e^{i n x}\right|^{2} d x=0
\end{aligned}
$$

Line 2 and Bessel's Inequality tell us that the series converges in $L^{2}(\mathbb{T})$. However, it doesn't converge pointwise everywhere on $\mathbb{T}$.

Ex. at $x=0, \sum_{n \neq 0} \frac{1}{|n|}$ diverges.

Proposition 7.7. Orthonormal Basis of $L^{2}(\mathbb{T})$
page 150

The Fourier basis elements are the functions

$$
e_{n}(x)=\frac{1}{\sqrt{2 \pi}} e^{i n x}
$$

$\left\{e_{n} \mid n \in \mathbb{Z}\right\}$ is an orthonormal basis of $L^{2}(\mathbb{T})$.

## Proof Outline

- Orthogonality

It is easily shown that

$$
\left\langle e_{m}, e_{n}\right\rangle= \begin{cases}1 & \text { if } m=n \\ 0 & \text { if } m \neq n\end{cases}
$$

- Completeness

This proof relies upon the ideas of convolution and approximate identities. (See Theorems 7.12 and 7.13.)

## Definition 7.8. Convolution

page 150

The convolution of two continuous functions $f, g: \mathbb{T} \rightarrow \mathbb{C}$ is the continuous function $f * g: \mathbb{T} \rightarrow \mathbb{C}$ defined by the integral

$$
(f * g)(x)=\int_{\mathbb{T}} f(x-y) g(y) d y
$$

Using the change of variable $z=x-y$, it is seen that

$$
(f * g)(x)=\int_{\mathbb{T}} f(z) g(x-z) d z
$$

so that $f * g=g * f$.

## Definition 7.9. Approximate Identity

Definition 7.1 on page 151

A family of functions $\left\{\varphi_{n} \in C(\mathbb{T}) \mid n \in \mathbb{N}\right\}$ is an approximate identity if
(a) $\varphi_{n}(x) \geq 0$
(b) $\int_{\mathbb{T}} \varphi_{n}(x) d x=1$
(c) $\lim _{n \rightarrow \infty} \int_{\delta \leq|x| \leq \pi} \varphi_{n}(x) d x=0 \quad \forall \delta>0$

Note: in (c), $\mathbb{T}$ is identified with $[-\pi, \pi]$.

## Theorem 7.10.

Theorem 7.2 on page 151 and Notes $1 / 5 / 11$ and FA 49

If $\left\{\varphi_{n} \in C(\mathbb{T}) \mid n \in \mathbb{N}\right\}$ is an approximate identity and $f \in C(\mathbb{T})$, then $\varphi_{n} * f$ converges uniformly to $f$ as $n \rightarrow \infty$.

Note: the term"approximate identity" comes from this result, since $\left\{\varphi_{n}\right\}$ is an approximation to the identity.

## Proof

$$
\begin{aligned}
f(x) & =\int_{\mathbb{T}} \varphi_{n}(y) f(x) d y \\
\left(\varphi_{n} * f\right)(x) & =\int_{\mathbb{T}} \varphi_{n}(y) f(x-y) d y \\
\left(\varphi_{n} * f\right)(x)-f(x) & =\int_{\mathbb{T}} \varphi_{n}(y)[f(x-y)-f(x)] d y
\end{aligned}
$$

- $f$ is uniformly continuous, so there exists $M$ such that $|f(x)| \leq M \forall x \in \mathbb{T}$
- $\exists \delta>0$ such that $|f(x)-f(y)| \leq \epsilon$ whenever $|x-y|<\delta$

$$
\begin{aligned}
\left|\left(\varphi_{n} * f\right)(x)-f(x)\right| & \leq \int_{-\pi}^{\pi} \varphi_{n}(y)|f(x-y)-f(x)| d y \\
& \leq \int_{|y|<\delta} \varphi_{n}(y)|f(x-y)-f(x)| d y+\int_{|y| \geq \delta} \varphi_{n}(y)|f(x-y)-f(x)| d y \\
& \leq \epsilon \int_{|y|<\delta} \varphi_{n}(y) d y+\int_{|y| \geq \delta} \varphi_{n}(y)[|f(x-y)|+|f(x)|] d y \\
& \leq \epsilon+2 M \int_{|y| \geq \delta} \varphi_{n}(y) d y
\end{aligned}
$$

Using property (c) of an approximate identity gives that $\varphi_{n} * f \rightarrow f$ uniformly in $C(\mathbb{T})$.

## Remark 7.11. Revised Approximate Identity Definition

Notes 1/5/11

More generally, $\varphi_{n} \in L^{1}(\mathbb{T})$ is an approximate identity if
(a) $\int_{\mathbb{T}}\left|\varphi_{n}(x)\right| d x \leq M \quad \forall n \in \mathbb{N}$
(b) $\int_{\mathbb{T}} \varphi_{n}(x) d x=1$
(c) $\lim _{n \rightarrow \infty} \int_{\delta \leq|x| \leq \pi} \varphi_{n}(x) d x=0 \quad \forall \delta>0$

Theorem 7.12. Weierstrass Approximation Theorem
Theorem 7.3 on page 152 and Notes $1 / 5 / 11$

The trigonometric polynomials are dense in $C(\mathbb{T})$ with respect to the uniform norm.

Proof

- Let $f \in C(\mathbb{T})$
- Generate an approximate identity that is a trigonometric polynomial
- Define $\varphi_{n}=c_{n}(1+\cos x)^{n}=c_{n}\left[2 \cos ^{2}\left(\frac{x}{2}\right)\right]^{n}$ and choose $c_{n}$ such that $\int_{\mathbb{T}} \varphi_{n}(x) d x=1$
- To show $\varphi_{n}$ is an approximate identity, we need to show that $\forall \delta>0, \lim _{n \rightarrow \infty} \int_{|x|>\delta} \varphi_{n}(x) d x=0$
* Fix $\epsilon>0 . \forall x, \delta \leq|x| \leq \pi, \exists r \in(0,1)$ such that

$$
\begin{aligned}
&(1+\cos x)<r(1+\cos y) \\
& \varphi_{n}(x)<r^{n} \varphi_{n}(y) \\
& \delta \varphi_{n}(x)<r^{n} \int_{-\delta / 2}^{\delta / 2} \varphi_{n}(y) d y \\
& \delta \varphi_{n}(x)<r^{n} \\
& 0 \leq \varphi_{n}(x)<\frac{r^{n}}{\delta} \quad \forall x \text { such that } \delta \leq|x| \leq \pi
\end{aligned}
$$

- So $\varphi_{n} \rightarrow 0$ uniformly on $\delta \leq|x| \leq \pi$ as $n \rightarrow \infty$, and $\int_{|x|>\delta} \varphi_{n}(x) d x \rightarrow 0$ as $n \rightarrow \infty$
- $\varphi_{n}$ is an approximate identity, so $\varphi_{n} * f$ is a trigonometric polynomial, and $\varphi_{n} * f$ converges uniformly to $f$ (See Theorem 7.10)

Corollary 7.13.
page 153 and 155 and Notes $1 / 5 / 11$

The trigonometric polynomials are dense in $L^{2}(\mathbb{T})$. That is, for any $f \in L^{2}(\mathbb{T})$,

$$
\begin{aligned}
& f(x)=\frac{1}{\sqrt{2 \pi}} \sum_{n=-\infty}^{\infty} \hat{f}_{n} e^{i n x} \\
& \hat{f}_{n}=\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{T}} f(x) e^{-i n x} d x
\end{aligned}
$$

If $f \in L^{2}(\mathbb{T})$ then the Fourier series of $f$ converges pointwise to $f$ a.e. (Carleson).

Proof
Let $f \in L^{2}(\mathbb{T})$.

- Choose $g \in C(\mathbb{T})$ such that $\|f-g\|_{L^{2}}<\epsilon / 2$. We can do this because $C(\mathbb{T})$ is dense in $L^{2}(\mathbb{T})$.
- Pick a trigonometric polynomial $p$ such that $\|g-p\|_{L^{2}}<\epsilon / 2 \sqrt{2 \pi}$.
- $\|g-p\|_{L^{2}}=\left(\int|g-p|^{2} d x\right)^{1 / 2} \leq\|g-p\|_{\infty} \sqrt{2 \pi}$
- $\|f-p\|_{L^{2}} \leq\|f-g\|_{L^{2}}+\|g-p\|_{L^{2}}<\epsilon / 2+\epsilon / 2$


## Corollary 7.14.

Notes 1/5/11
$\left\{\left.\frac{1}{\sqrt{2 \pi}} e^{i n x} \right\rvert\, n \in \mathbb{Z}\right\}$ is an orthonormal basis of $L^{2}(\mathbb{T})$.

Definition 7.15. Periodic Fourier Transform page 153 and Notes $1 / 7 / 11$

The Periodic Fourier Transform $\mathcal{F}: L^{2}(\mathbb{T}) \rightarrow \ell^{2}(\mathbb{Z})$ maps a function to its sequence of Fourier coefficients by

$$
\mathcal{F} f=\left(\hat{f}_{n}\right)_{n=-\infty}^{\infty}
$$

Thus, the $L^{2}$ norm of a function can be computed by

$$
\int_{\mathbb{T}}|f(x)|^{2} d x=\sum_{n=-\infty}^{\infty}\left|\hat{f}_{n}\right|^{2}
$$

This implies that $\left(\hat{f}_{n}\right) \in \ell^{2}(\mathbb{Z})$. Furthermore, the Projection Theorem ( 6.13 in the book) implies that

$$
f_{N}(x)=\frac{1}{\sqrt{2 \pi}} \sum_{n=-N}^{N} \hat{f}_{n} e^{i n x}
$$

is the best approximation of $f$ by a trigonometric polynomial of degree $N$ in the $L^{2}$-norm.

Theorem 7.16. Parseval's Theorem
Notes 1/7/11

Given $f, g \in L^{2}(\mathbb{T})$, then

$$
\begin{aligned}
f(x) & =\frac{1}{\sqrt{2 \pi}} \sum_{n \in \mathbb{Z}} a_{n} e^{i n x} \\
g(x) & =\frac{1}{\sqrt{2 \pi}} \sum_{n \in \mathbb{Z}} b_{n} e^{i n x} \\
\langle f, g\rangle & =\sum_{n \in \mathbb{Z}} \overline{a_{n}} b_{n}
\end{aligned}
$$

## Proposition 7.17.

Proposition 7.4 on page 154

If $f, g \in L^{2}(\mathbb{T})$, then $f * g \in C(\mathbb{T})$ and

$$
\|f * g\|_{\infty} \leq\|f\|_{2}\|g\|_{2}
$$

Proof

$$
(f * g)(x)=\int_{\mathbb{T}} f(x-y) g(y) d y
$$

If $f, g \in C(\mathbb{T})$, then we can apply the Cauchy-Schwarz Inequality to get

$$
|f * g(x)| \leq\|f\|_{L^{2}}\|g\|_{L^{2}}
$$

Taking the supremum of both sides yields

$$
\|f * g\|_{\infty} \leq\|f\|_{L^{2}}\|g\|_{L^{2}}
$$

If $f, g \in L^{2}(\mathbb{T})$, then there exist sequences $\left(f_{k}\right),\left(g_{k}\right) \in C(\mathbb{T})$ such that $\left\|f-f_{k}\right\|_{2} \rightarrow 0$ and $\left\|g-g_{k}\right\|_{2} \rightarrow 0$ as $k \rightarrow \infty$. Also, the sequence $\left(f_{k} * g_{k}\right) \in C(\mathbb{T})$ is Cauchy with respect to the sup-norm, since

$$
\begin{aligned}
\left\|f_{j} * g_{j}-f_{k} * g_{k}\right\| & \leq\left\|\left(f_{j}-f_{k}\right) * g_{j}\right\|_{\infty}+\left\|f_{k} *\left(g_{j}-g_{k}\right)\right\|_{\infty} \\
& \leq\left\|f_{j}-f_{k}\right\|_{2}\left\|g_{j}\right\|_{2}+\left\|f_{k}\right\|_{2}\left\|g_{j}-g_{k}\right\|_{2} \\
& \leq M\left(\left\|f_{j}-f_{k}\right\|_{2}+\left\|g_{j}-g_{k}\right\|_{2}\right)
\end{aligned}
$$

where $M \geq\left\|f_{j}\right\|_{2}$ and $M \geq\left\|g_{k}\right\|_{2}$, since the sequences converge in $L^{2}(\mathbb{T})$. Since $C(\mathbb{T})$ is complete, the sequence $\left(f_{k} * g_{k}\right)$ converges uniformly to a continuous function $f * g \in C(\mathbb{T})$, and $f * g$ satisfies the inequality.

Theorem 7.18. Convolution Theorem
Theorem 7.5 on page 154 and Notes $1 / 10 / 11$

If $f, g \in L^{2}(\mathbb{T})$, then

$$
\begin{array}{ll}
\text { (Book) } & \widehat{(f * g)}_{n}=\sqrt{2 \pi} \hat{f}_{n} \hat{g}_{n} \\
\text { (Notes) } & \widehat{(f * g)}_{n}=\hat{f}_{n} \hat{g}_{n}
\end{array}
$$

## Proof Outline

Compute $\widehat{(f * g)}{ }_{n}$, using Fubini's Theorem to change the order of integration.

Remark 7.19. Alternative bases for $L^{2}$
page 155 and Notes $1 / 7 / 11$

The non-normalized orthogonal basis:

$$
\begin{gathered}
\left\{e^{i n x}\right\} \\
\hat{f}_{n}=\frac{1}{2 \pi} \int_{\mathbb{T}} f(x) e^{-i n x} d x
\end{gathered}
$$

Sines and Cosines:

$$
\begin{gathered}
\{1, \cos (n x), \sin (n x) \mid n=1,2,3, \ldots\} \\
f(x)=\frac{1}{2} a_{0}+\sum_{n=1}^{\infty}\left(a_{n} \cos (n x)+b_{n} \sin (n x)\right) \\
a_{0}=\frac{1}{\pi} \int_{\mathbb{T}} f(x) d x \quad a_{n}=\frac{1}{\pi} \int_{\mathbb{T}} f(x) \cos (n x) d x \quad b_{n}=\frac{1}{\pi} \int_{\mathbb{T}} f(x) \sin (n x) d x
\end{gathered}
$$

## 7.2 $\quad L^{1}$ Functions

## Remark 7.20. L ${ }^{1}$ Functions

Notes $1 / 7 / 11$
$L^{1}(\mathbb{T})$ is the space of periodic functions $f: \mathbb{T} \rightarrow \mathbb{C}$ such that

$$
\|f\|_{L^{1}}=\int_{\mathbb{T}}|f(x)| d x<\infty
$$

Note that $L^{1}(\mathbb{T})$ is a Banach space but not a Hilbert space. We can define the Fourier coefficients of $f$ as

$$
c_{n}=\int_{\mathbb{T}} f(x) e^{-i n x} d x
$$

Note that $\left|c_{n}\right| \leq \int|f(x)| d x$. We can write the trigonometric polynomial approximation of $f$ as

$$
f(x) \sim \sum_{n \in \mathbb{Z}} c_{n} e^{i n x}
$$

However, this does not necessarily converge to $f$.

## Lemma 7.21. Riemann-Lebesgue Lemma

Notes $1 / 7 / 11$ and $1 / 10 / 11$

If $f \in L^{1}(\mathbb{T})$ has Fourier coefficients $c_{n}$, then $c_{n} \rightarrow 0$ as $|n| \rightarrow \infty$.
Proof Outline ( $1 / 7 / 11$ )

- Prove for smooth functions (use Integration By Parts)
- Approximate non-smooth functions with smooth functions
$\underline{\text { Proof Outline (1/10/11) }}$
- Fix $\epsilon>0$
- The trigonometric polynomials are dense in $L^{1}(\mathbb{T})$, so we can pick a trigonometric polynomial $p$ such that $\|f-p\|_{L^{1}}<\epsilon$
- If $\operatorname{deg} p=N$ and $n>N$, then

$$
\begin{aligned}
|\hat{f}(n)| & =\frac{1}{2 \pi}\left|\int f e^{-i n x} d x\right| \\
& =\frac{1}{2 \pi}\left|\int(f-p) e^{-i n x} d x\right| \quad \text { Note: } \int p e^{-i n x} d x=0 \forall n>N \text { by orthogonality } \\
& \leq \frac{1}{2 \pi}\|f-p\|_{L^{1}} \\
& \leq \frac{\epsilon}{2 \pi}<\epsilon
\end{aligned}
$$

Definition 7.22. Fourier Transform for $L^{1}(\mathbb{T})$
Notes 1/10/11

The Fourier Transform $\mathcal{F}: f \rightarrow \hat{f}, \mathcal{F}: L^{1}(\mathbb{T}) \rightarrow C_{0}(\mathbb{Z})$

$$
\begin{aligned}
& C_{0}(\mathbb{Z})=\left\{\left(c_{n}\right)_{n \in \mathbb{Z}} \mid c_{n} \rightarrow 0 \text { as }|n| \rightarrow \infty\right\} \\
& \left\|\left(c_{n}\right)\right\|_{\infty}=\max _{n \in \mathbb{Z}}\left|c_{n}\right|
\end{aligned}
$$

$\mathcal{F}$ is a bounded linear map, with $\|\mathcal{F} f\|_{\infty} \leq\|f\|_{L^{1}}$ Note: $\mathcal{F}$ is not onto.

Example 7.23. $\mathcal{F}$ is not onto
Notes 1/10/11

There is no function whose Fourier coefficients are

$$
\hat{f}(n)=\frac{i \operatorname{sgn}(n)}{\log |n|} \quad|n| \geq 2
$$

### 7.3 Kernels and Summability Methods

## Definition 7.24. Dirichlet Kernel

Notes 1/10/11 and FA 44

The Dirichlet kernel is

$$
\begin{aligned}
& D_{N}(x)=\frac{1}{2 \pi} \sum_{|n| \leq N} e^{i n x}=\frac{1}{2 \pi}\left[\frac{\sin \left(\left(N+\frac{1}{2}\right) x\right)}{\sin \left(\frac{x}{2}\right)}\right] \quad x \neq 0 \\
& D_{N}(0)=\frac{1}{2 \pi}(2 N+1)
\end{aligned}
$$

(See the Kernel Overview.)
Derivation of the Dirichlet Kernel

Suppose $f \in L^{1}(\mathbb{T}), f(x) \sim \sum \hat{f}_{n} e^{i n x}$. Define the $N^{\text {th }}$ partial sum of the Fourier series of $f$ as

$$
\begin{aligned}
S_{N}(f)(x) & =\sum_{|n| \leq N} \hat{f}_{n} e^{i n x} \\
& =\frac{1}{2 \pi} \sum_{|n| \leq N}\left(\int f(y) e^{-i n y} d y\right) e^{i n x} \\
& =\frac{1}{2 \pi} \int\left(\sum_{|n| \leq N} e^{i n(x-y)}\right) f(y) d y \\
& =\int D_{N}(x-y) f(y) d y=D_{N} * f
\end{aligned}
$$



Figure 1: Dirichlet kernels.

Example 7.25. $D_{N}$ is not an approximate identity
Notes 1/12/11

The Dirichlet kernel is not an approximate identity.

$$
\begin{aligned}
& \text { (a) } \int D_{N} d x=\int\left(\frac{1}{2 \pi} \sum e^{i n x}\right) d x=\frac{1}{2 \pi} \cdot 2 \pi=1 \\
& \text { (b) } \int \frac{4}{\pi^{2}} \sum_{k=1}^{N} \frac{1}{k} \leq\left|D_{N}\right| d x \leq \frac{4}{\pi^{2}}\left(\sum_{k=1}^{N} \frac{1}{k}\right)+2+\frac{\pi}{4} \\
& \text { As } N \rightarrow \infty, \int\left|D_{N}\right| d x=\frac{4}{\pi} \log N+O(1) \rightarrow \infty \text { as } N \rightarrow \infty \\
& \text { (c) } \text { For } \delta>0, \lim _{N \rightarrow \infty} \int_{|x|>\delta}\left|D_{N}\right| d x \nrightarrow 0
\end{aligned}
$$

Thus, we can't conclude that if $f \in C(\mathbb{T})$ or $f \in L^{1}(\mathbb{T})$ then $D_{N} * f \rightarrow f$ uniformly

Theorem 7.26. Absolute Convergence
HW 3 Problem 2 and FA page 41

If $f \in C(\mathbb{T})$ and its Fourier series is absolutely convergent, $\sum_{n \in \mathbb{Z}}|\hat{f}(n)|<\infty$, then the Fourier series converges uniformly to $f$.

Let $\mathcal{A}(\mathbb{T})$ denote the space of integrable functions whose Fourier coefficients are absolutely convergent. That is, $f \in \mathcal{F}(\mathbb{T})$ if $\sum_{n \in \mathbb{Z}}|\hat{f}(n)|<\infty$. If $f \in \mathcal{A}(\mathbb{T})$, then $f \in C(\mathbb{T})$.

## Definition 7.27. Summability Method: Cesáro Summation

Notes $1 / 12 / 11$ and FA 52

The $N^{\text {th }}$ Cesáro sum of a series is the average of the first $N$ partial sums in the series:

$$
\sigma_{N}=\frac{s_{0}+s_{1}+\ldots+s_{N-1}}{N}
$$

Example 7.28. Cesáro Summation Example
Notes 1/12/11

Consider the series $\sum_{n=1}(-1)^{n}=1-1+1-1+1 \ldots$. Then the $n$th partial sum is

$$
S_{N}= \begin{cases}1 & N \text { odd } \\ 0 & N \text { even }\end{cases}
$$

Consider the averages of partial sums:

$$
\begin{aligned}
& \sigma_{N}=\frac{S_{1}+\ldots+S_{N}}{N} \\
& \sigma_{N}=\left\{\begin{array}{cl}
\frac{1}{\frac{1}{2}(N+1)} \\
N & \frac{1}{2}+\frac{1}{2 N}
\end{array} \quad N\right. \text { even }
\end{aligned} \rightarrow \frac{1}{2} \text { as } N \rightarrow \infty
$$

Thus, $\sum_{n=1}(-1)^{n}=\frac{1}{2}(\mathrm{C})$.

## Theorem 7.29.

Notes 1/14/11

Cesáro summation is regular, meaning that if $\sum a_{n}=s$ then $\sum a_{n}=s$ (C).

Definition 7.30. Fejér Kernel
Notes 1/12/11

The Fejér Kernel is:

$$
\begin{aligned}
& K_{N}(x)=\frac{1}{2 \pi} \sum_{|n| \leq N}\left(1-\frac{|n|}{N+1}\right) e^{i n x} \\
& K_{N}(x)=\frac{1}{2 \pi(N+1)}\left[\frac{\sin \left(\frac{(N+1) x}{2}\right)}{\sin \left(\frac{x}{2}\right)}\right]^{2}
\end{aligned}
$$

(See the Kernel Overview.)
Proof (that the two forms are equivalent)

- Consider

$$
\left[\frac{1}{2}\left(e^{i x}+e^{-i x}\right)-1\right] K_{N}(x)=\frac{1}{2 \pi N}\left(\frac{1}{2} e^{i(N+1) x}+\frac{1}{2} e^{-i(N+1) x}-1\right)
$$

- Use the fact that

$$
\left(\sin \frac{x}{2}\right)^{2}=-\frac{1}{4}\left(e^{i x}-2+e^{-i x}\right)
$$

Derivation of the Fejér Kernel
Form the $N^{\text {th }}$ Cesáro mean of the Fourier series:

$$
\begin{aligned}
\sigma_{N}(f)(x) & =\frac{S_{0} f+S_{1} f+\ldots+S_{N} f}{N+1} \\
& =\frac{1}{2 \pi} \sum_{|n| \leq N}\left(1-\frac{|n|}{N+1}\right) \hat{f}(n) e^{i n x} \\
& =K_{N} * f
\end{aligned}
$$



Figure 2: Fejér kernels.

## Theorem 7.31.

Notes 1/12/11
$K_{N}$ is an approximate identity. If $f \in C(\mathbb{T})$, then $\sigma_{N} f=K_{N} * f \rightarrow f$ uniformly and if $f \in L^{p}(\mathbb{T})$, then $\sigma_{N} f=K_{N} * f \rightarrow f$ in $L^{p}(\mathbb{T})$.

Corollary 7.32. $1 / 12 / 11$

Suppose $f, g \in L^{1}(\mathbb{T})$ and $\hat{f}=\hat{g}$. Then $f=g$.

Proof

- Set $h=f-g$
- Then $\hat{h}(n)=0$
- $K_{N} * h \rightarrow h$ in $L^{1}$
- $K_{N} * h=0 \forall N$, so $h=0 \Rightarrow f=g$

Note: we could have used the original approximate identity for this proof.
Definition 7.33. Summability Method: Abel Summation
Notes $1 / 14 / 11$

$$
\begin{align*}
& S=\sum_{n=0}^{\infty} a_{n} \\
& S=\lim _{r \rightarrow 1^{-}} \sum_{n=0}^{\infty} a_{n} r^{n} \tag{A}
\end{align*}
$$

Theorem 7.34.
Notes $1 / 14 / 11$

Abel summation is regular.

Proof

- We will use summation by parts. Suppose $S=\sum_{n=0}^{\infty} a_{n}, S_{n}=\sum_{k=0}^{n} a_{k}, S_{n} \rightarrow S$ as $n \rightarrow \infty$

$$
\begin{array}{rlrl}
\sum_{n=0}^{\infty} a_{n} r^{n} & =a_{0}+\sum_{n=1}^{\infty}\left(S_{n}-S_{n-1}\right) r^{n} & \left(\text { Since } a_{n}=S_{n}-S_{n-1}\right) \\
& =a_{0}+\sum_{n=1}^{\infty}\left(S_{n}-S_{n} r^{n+1}\right) & \quad \text { (re-index) } \\
& =a_{0}+(1-r) \sum_{n=1}^{\infty}\left(S_{n} r^{n}\right)-S_{0} r & \\
& =(1-r) \sum_{n=0}^{\infty} S_{n} r^{n} & \left(S_{0}=a_{0}\right) \\
\left|\sum_{n=0}^{\infty}\left(a_{n} r^{n}\right)-s\right| & =(1-r)\left|\sum_{n=0}^{\infty}\left(S_{n}-S\right) r^{n}\right| \leq(1-r) \sum_{n=0}^{\infty}\left|S_{n}-S\right| r^{n} & S=(1-r) \sum_{n=0}^{\infty} r^{n} \\
& & S-r) \sum_{n=0}^{\infty} S r^{n}
\end{array}
$$

- Fix $\epsilon>0$. Choose $N$ such that $\left|S_{n}-S\right|<\epsilon / 2$ for $n>N$. Then

$$
\left|\sum_{n=0}^{\infty} a_{n} r^{n}-S\right|<(1-r) \sum_{n=0}^{N}\left|S_{n}-S\right| r^{n}+\frac{\epsilon}{2} \underbrace{(1-r) \sum_{n=N+1}^{\infty} r^{n}}_{\leq 1}
$$

- Choose $(1-r)<\delta$, where $\delta \sum_{n=0}^{N}\left|S_{n}-S\right|<\epsilon / 2$
- $n>N \Rightarrow\left|\sum_{n=0}^{\infty} a_{n} r^{n}-S\right|<\epsilon / 2+\epsilon / 2=\epsilon$

Theorem 7.35. Tauber $\mathcal{E}$ Littlwood
Notes 1/14/11

Suppose that $\lim _{r \rightarrow 1^{-}} \sum_{n=0}^{\infty} a_{n} r^{n}$ exists and $n a_{n}=O(1)$ as $n \rightarrow \infty$. (i.e. there is an $M$ such that $\left|n a_{n}\right| \leq M \forall n$.) Then $\sum a_{n}$ exists (and is equal to the limit).

Definition 7.36. Poisson Kernel
Notes $1 / 14 / 11$

Identify $\mathbb{T}$ as the unit circle in $\mathbb{C}$, i.e.

$$
\begin{aligned}
& \mathbb{T}=\{z \in \mathbb{C}| | z \mid=1\} \Leftrightarrow z=e^{i \theta} \\
& f(\theta) \sim \sum_{n \in \mathbb{Z}} \hat{f}(n) e^{i n \theta} \\
& \hat{f}(n)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f(\theta) e^{-i n \theta} d \theta \\
& f_{r}(\theta)=\sum_{n \in \mathbb{Z}} \hat{f}(n) r^{|n|} e^{i n \theta} \\
&=P_{r} * f(\theta)
\end{aligned}
$$

The Poisson kernel is

$$
\begin{aligned}
& P_{r}(\theta)=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{i n \theta}, \quad 0<r<1 \\
& P_{r}(\theta)=\frac{1}{2 \pi}\left[\frac{1-r^{2}}{1-2 r \cos \theta+r^{2}}\right] \\
& P_{r}(0)=\frac{1}{2 \pi} \frac{1-r^{2}}{(1-r)^{2}}
\end{aligned}
$$

(See the Kernel Overview.)


Figure 3: Poisson kernels.

Remark 7.37. Properties of the Poisson Kernel
Notes 1/14/11

- The Poisson kernel is not a trigonometric polynomial
- The Poisson kernel satisfies:
(a) $\int P_{r}(\theta) d \theta=1$
(b) $P_{r} \geq 0$
(c) $P_{r}(\theta) \rightarrow 0$ uniformly as $r \rightarrow 1^{-}$on $\delta<|\theta|<\pi$


## Theorem 7.38.

Notes 1/14/11
$P_{r}$ is an approximate identity as $r \rightarrow 1^{-}$.

Corollary 7.39.
Notes 1/14/11

If $f \in L^{p}(\mathbb{T}), 1 \leq p<\infty$, then $P_{r} * f \rightarrow f$ as $r \rightarrow 1^{-}$.
If $f \in C(\mathbb{T})$, then $P_{r} * f \rightarrow f$ uniformly.

## Remark 7.40. Kernel Overview

## Dirichlet

- Equations:

$$
\begin{aligned}
& -D_{N}(x)=\frac{1}{2 \pi} \sum_{|n| \leq N} e^{i n x} \\
& -D_{N}(x)=\frac{1}{2 \pi}\left[\frac{\sin \left(\left(N+\frac{1}{2}\right) x\right)}{\sin \left(\frac{x}{2}\right)}\right], \quad x \neq 0 \\
& -D_{N}(0)=\frac{1}{2 \pi}(2 N+1)
\end{aligned}
$$

- Summability Method: Standard
- Approximate Identity: No


## Fejér

- Equations:

$$
\begin{aligned}
& -K_{N}(x)=\frac{1}{2 \pi} \sum_{|n| \leq N}\left(1-\frac{|n|}{N+1}\right) e^{i n x} \\
& -K_{N}(x)=\frac{1}{2 \pi(N+1)}\left[\frac{\sin \left(\frac{(N+1) x}{2}\right)}{\sin \left(\frac{x}{2}\right)}\right]^{2}
\end{aligned}
$$

- Summability Method: Cesáro
- Approximate Identity: Yes


## Poisson

- Equations:
$-P_{r}(\theta)=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{i n \theta}, \quad 0<r<1$
$-P_{r}(\theta)=\frac{1}{2 \pi}\left[\frac{1-r^{2}}{1-2 r \cos \theta+r^{2}}\right]$
$-P_{r}(0)=\frac{1}{2 \pi} \frac{1-r^{2}}{(1-r)^{2}}$
- Summability Method: Abel
- Approximate Identity: Yes, as $r \rightarrow 1^{-}$


### 7.4 Harmonic Functions

## Definition 7.41. Harmonic

Notes 1/19/11

Let $\Omega \subset \mathbb{R}^{n}$ be an open set.
$u: \Omega \rightarrow \mathbb{R}$ is harmonic on $\Omega$ if $\Delta u=0$ in $\Omega$.
Recall: $\Delta=\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\ldots+\frac{\partial^{2}}{\partial x_{n}^{2}}$

## Remark 7.42. Harmonic ${ }^{6}$ Analytic Functions

Notes 1/19/11

There is a close connection in 2-D between harmonic and analytic (holomorphic) functions.

$$
\begin{aligned}
& F: \Omega \rightarrow \mathbb{C} \\
& F(z)=u(x, y)+i v(x, y)
\end{aligned}
$$

where $u, v$ satisfy the Cauchy-Riemann equations:

$$
\left.\begin{array}{c}
u_{x}=v_{y} \\
u_{y}=-v_{x}
\end{array}\right\} \Rightarrow u_{x x}+v_{y y}=0
$$

## Example 7.43. $\Delta u=0$ on the Complex Unit Disk

Notes 1/19/11

Consider the Dirichlet problem on $D=\left\{(x, y) \subset \mathbb{R}^{2} \mid x^{2}+y^{2}<1\right\}$ :

$$
\begin{aligned}
& \Delta u=0 \text { in } D \\
& u=f \text { on } \partial D=\pi
\end{aligned}
$$

Here $f \in C(\partial D)$.
Want $u \in C^{2}(D) \cap C(\bar{D})$.
Use separation of variables:

$$
u(r, \theta)=F(r) G(\theta)
$$

We get that:

$$
\begin{aligned}
& G(\theta)=e^{i n \theta} \\
& F(r)=A r^{n}+B r^{-n} \quad n \neq 0 \\
& F(r)=A+B \ln r \quad n=0
\end{aligned}
$$

We want the solution to belong to $C^{2}(D)$, so we set

$$
\begin{aligned}
F(r) & =r^{|n|}, \quad n \in \mathbb{Z} \\
\Rightarrow u(r, \theta) & =\sum_{n \in \mathbb{Z}} c_{n} r^{|n|} e^{i n \theta}
\end{aligned}
$$

We want that:

$$
\begin{aligned}
u(1, \theta) & =f(\theta)=\sum_{n \in \mathbb{Z}} c_{n} e^{i n \theta} \\
\Rightarrow c_{n} & =\hat{f}(n)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f(\theta) e^{-i n \theta} d \theta
\end{aligned}
$$

Note that:

$$
u(r, \theta)=\underbrace{\left(P_{r} * f\right)(\theta)}_{\text {Green's function }}=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{i n \theta}
$$

## Remark 7.44.

Notes 1/19/11
$P_{r}(\theta)$ is a $C^{\infty}(D)$ function of $r, \theta$ in $0 \leq r<1$, and

$$
\Delta P_{r}(\theta)=\frac{1}{r} \frac{\partial}{\partial r}\left(r \frac{\partial P_{r}}{\partial r}\right)+\frac{1}{r^{2}} \frac{\partial^{2} P_{r}}{\partial \theta^{2}}=0
$$

## Theorem 7.45.

Notes 1/19/11

Suppose that $f \in C(\partial D)$. Then $u(r, \theta)=\left(P_{r} * f\right)(\theta)$ is a solution of

$$
\left\{\begin{array}{cl}
\Delta u=0 & \text { in } D \\
u=f & \text { on } \partial D
\end{array}\right.
$$

Moreover, $u \in C^{\infty}(D) \cap C(\bar{D})$.
Proof

- $u(r, \theta)=\int_{\mathbb{T}} P_{r}(\theta-\phi) f(\phi) d \phi$ (by Lebesgue Dominated Convergence Theorem)
- So $u \in C^{\infty}(D)$, and $\Delta u=0$
- Moreover, $P_{r} * f \rightarrow f$ uniformly as $r \rightarrow 1^{-}$
- So $u \in C(\bar{D})$


## Theorem 7.46.

Notes 1/19/11

There is a unique solution $u \in C^{2}(D) \cap C(\bar{D})$ of the Dirichlet problem. (Can be proved using the maximum principle and/or energy estimates.)

## Corollary 7.47.

Notes 1/19/11

Every harmonic function $u \in C^{2}(D) \cap C(\bar{D})$ is smooth and has the mean value property:

$$
u(r=0)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f(\theta) d \theta
$$

### 7.5 Hausdorff-Young Inequality

## Remark 7.48. Background Info/Review

Notes 1/21/11

- Function Spaces
- Let $1 \leq p<\infty$. If $f \in L^{p}(\mathbb{T})$, then $f: \mathbb{T} \rightarrow \mathbb{C}$ and $\|f\|_{p}=\left(\int_{\mathbb{T}}|f|^{p} d x\right)^{1 / p}<\infty$.
$-f=g$ in $L^{p}$ if $f=g$ a.e.
- In $L^{\infty},\|f\|_{\infty}=\operatorname{ess} \sup _{\mathbb{T}}|f(x)|=\inf _{\text {measure } N=0} \sup \{|f(x)| \mid x \in \mathbb{T} \backslash N\}$
- Sequence Spaces
- Let $1 \leq q<\infty$. If $\hat{f} \in \ell^{q}(\mathbb{Z}), \hat{f}: \mathbb{Z} \rightarrow \mathbb{C}$, then $\|\hat{f}\|_{q}=\left(\sum_{n \in \mathbb{Z}}|\hat{f}(n)|^{q}\right)^{1 / q}<\infty$
$-\operatorname{In} \ell^{\infty},\|\hat{f}\|_{\infty}=\sup _{n \in \mathbb{Z}}|\hat{f}(n)|$
- Question: When is $\mathcal{F}: L^{p}(\mathbb{T}) \rightarrow \ell^{q}(\mathbb{Z}), f \mapsto \hat{f}$, a bounded linear map?
$-\mathcal{F}: L^{2} \rightarrow \ell^{2}$
* $\|\mathcal{F} f\|_{\ell^{2}}=\frac{1}{\sqrt{2 \pi}}\|f\|_{L^{2}}$
* $\mathcal{F}$ is onto
$-\mathcal{F}: L^{1} \rightarrow C_{0} \subset \ell^{\infty}$
* $\|\mathcal{F} f\|_{\ell^{\infty}} \leq \frac{1}{2 \pi}\|f\|_{L^{1}}$

Theorem 7.49. Hausdorff-Young Theorem/Inequality
Notes 1/21/11

Suppose $1 \leq p \leq 2$ and $2 \leq p^{\prime} \leq \infty$ are Hölder conjugates $\left(\frac{1}{p}+\frac{1}{p^{\prime}}=1\right)$.
Then $\mathcal{F}: L^{p}(\mathbb{T}) \rightarrow \ell^{p^{\prime}}(\mathbb{Z})$ is a bounded linear map, i.e. $\|\hat{f}\|_{\ell^{p^{\prime}}} \leq C_{p}\|f\|_{L^{p}}$.

## Remark 7.50.

Notes 1/21/11

1. Interpolation result (Riesz-Thorin Theorem)
2. $\mathcal{F}$ is not onto if $1 \leq p<2$.

- Ex: $p=1, p^{\prime}=\infty$, then $f \in L^{1} \rightarrow \hat{f} \in C_{0} \Rightarrow$ not all of $\ell^{\infty}$
- $\sum_{|n| \geq 2} \frac{i \operatorname{sgn}(n)}{\log n} e^{i n x}$ is not the Fourier series of any $L^{1}$ function

3. This result does not hold for $2<p \leq \infty$
4. If $f \in L^{p}$ (or even if $f \in C$ ), one can't say much about the Fourier coefficients $\hat{f}$ beyond the fact that $f \in L^{p}$ so $\hat{f} \in \ell^{2}$

- Example:

$$
\begin{aligned}
f(x) & =\sum_{n=2}^{\infty} \frac{e^{i n \log n}}{n^{1 / 2}(\log n)^{2}} e^{i n x} \\
\hat{f}(n) & =\frac{e^{i n \log n}}{n^{1 / 2}(\log n)^{2}} \\
\sum|\hat{f}(n)|^{2} & =\sum \frac{1}{n(\log n)^{4}}<\infty
\end{aligned}
$$

$\hat{f} \in \ell^{2}$ so $f \in L^{2}$. Is $\hat{f} \in \ell^{p}$ for $p<2$, e.g. $p=2-\epsilon$ ?

$$
\sum|\hat{f}(n)|^{2-\epsilon}=\sum \frac{1}{n^{1-\epsilon / 2}(\log n)^{4-2 \epsilon}}=\infty
$$

So $\hat{f} \notin \ell^{p^{\prime}}$ for any $p^{\prime}<2$

### 7.6 Fourier Series of Differentiable Functions (Section 7.2 in H\&N)

Definition 7.51. Fourier Series Differentiation
Notes 1/24/11

$$
\begin{aligned}
f(x) & =\sum_{n \in \mathbb{Z}} c_{n} e^{i n x} \\
f^{\prime}(x) & =\sum_{n \in \mathbb{Z}} i n c_{n} e^{i n x} \\
\mathcal{F} & : \frac{d}{d x} \mapsto i n
\end{aligned}
$$

## Proposition 7.52.

Notes 1/24/11

If $f \in C^{1}(\mathbb{T})$, then

$$
\widehat{f}^{\prime}(n)=i n \hat{f}(n)
$$

(Actually, it is sufficient that $f \in L^{1}(\mathbb{T})$.)

See Definition 7.56 and Proposition 11.21.

## Definition 7.53. Orders

Notes 1/24/11

If $\phi, \psi: \mathbb{Z} \rightarrow \mathbb{C}$, we say that

- $\phi=O(\psi)$ as $|n| \rightarrow \infty$ if there exists $C$ such that $|\phi(n)| \leq C|\psi(n)| \forall n \in \mathbb{Z}$
- $\phi=o(\psi)$ as $|n| \rightarrow \infty$ if $\lim _{|n| \rightarrow \infty}\left|\frac{\phi(n)}{\psi(n)}\right|=0$


## Theorem 7.54.

Notes 1/24/11

If $f \in C^{1}(\mathbb{T})$, then $\hat{f}(n)=o\left(\frac{1}{n}\right)$ as $|n| \rightarrow \infty$
If $f \in C^{k}(\mathbb{T})$, where $k \in \mathbb{N}$, then $\hat{f}(n)=o\left(\frac{1}{n^{k}}\right)$ as $|n| \rightarrow \infty$

## Proof

- $\widehat{f}^{\prime}(n)=i n \hat{f}(n)$ if $f \in C^{1}$
- $\hat{f}(n)=\frac{1}{i n} \widehat{f}^{\prime}(n), n \neq 0$, and $\widehat{f}^{\prime}(n) \rightarrow 0$ as $|n| \rightarrow \infty$ by the Riemann-Lebesgue Lemma
- So $\hat{f}(n)=o\left(\frac{1}{n}\right)$ as $|n| \rightarrow \infty$
- In general, $\hat{f}(n)=\frac{1}{(i n)^{k}} \widehat{f^{k}}(n)=o\left(\frac{1}{n^{k}}\right)$


## Corollary 7.55.

page 157 and Notes $1 / 24 / 11$

If $f \in C^{\infty}(\mathbb{T})$, then $\lim _{|n| \rightarrow \infty}|n|^{k} \hat{f}(n)=0 \forall k \in \mathbb{N}$.
In other words, the Fourier coefficients of smooth functions form a rapidly decreasing sequence that decreases faster than any polynomial. Heuristically, a smooth function contains a small amount of high frequency components.

Compare to Theorem 11.18.

Definition 7.56. Weak $L^{2}$-derivatives (1)
Notes 1/24/11

Suppose that $f \in L^{2}(\mathbb{T})$ such that $\sum_{n \in \mathbb{Z}} n^{2}|\hat{f}(n)|^{2}<\infty$. Then we define the weak $L^{2}$-derivative $g=f^{\prime} \in L^{2}(\mathbb{T})$ by

$$
g(x)=\sum_{n \in \mathbb{Z}} i n \hat{f}(n) e^{i n x}
$$

See Proposition 7.52 and Proposition 11.21.

Definition 7.57. Sobolev Space (1)
page 158 and Notes $1 / 24 / 11$

$$
\begin{aligned}
& H^{1}(\mathbb{T})=\left\{f \in L^{2}(\mathbb{T}) \mid f^{\prime} \in L^{2}(\mathbb{T})\right\} \\
& \langle f, g\rangle_{H^{1}}=\int_{\mathbb{T}}\left(\bar{f} g+\overline{f^{\prime}} g^{\prime}\right) d x=\sum_{n \in \mathbb{Z}}\left(1+n^{2}\right) \overline{\hat{f}(n)} g(n) \\
& \|f\|_{H^{1}}=\left[\int_{\mathbb{T}}\left(|f|^{2}+\left|f^{\prime}\right|^{2}\right) d x\right]^{1 / 2}
\end{aligned}
$$

In other words, $f \in H^{1}(\mathbb{T})$ iff $f$ and its weak derivative $f^{\prime}$ (defined by integration by parts) belong to $L^{2}(\mathbb{T})$.

## Definition 7.58. Integration By Parts

Notes $1 / 24 / 11$

For $f, g \in H^{1}$ :

$$
\begin{aligned}
\int_{\mathbb{T}} \overline{f^{\prime}} g d x & =2 \pi \sum \overline{\hat{f}^{\prime}(n)} \hat{g}(n) \\
& =2 \pi \sum \overline{\operatorname{in\hat {f}}(n)} \hat{g}(n) \\
& =-2 \pi \sum \overline{\hat{f}(n)} i n \hat{g}(n) \\
& =-2 \pi \sum \overline{\hat{f}(n)} \widehat{g^{\prime}}(n) \\
& =-\int_{\mathbb{T}} \bar{f} g^{\prime} d x
\end{aligned}
$$

Definition 7.59. Weak Derivative (2) page 159 and Notes $1 / 24 / 11$

A function $g \in L^{1}(\mathbb{T})$ is the weak derivative of a function $f \in L^{1}(\mathbb{T})$, written $g=f^{\prime}$, if for every $\phi \in C^{\infty}(\mathbb{T})$ we have

$$
\int_{\mathbb{T}} f \phi^{\prime} d x=-\int_{\mathbb{T}} g \phi d x
$$

In other words, we are using integration by parts $\left(\int_{\mathbb{T}} \overline{f^{\prime}} g d x=-\int_{\mathbb{T}} \bar{f} g^{\prime} d x\right)$, to define $f^{\prime}$ pointwise a.e. We determine $\hat{g}(n) \forall n$ by choosing $\phi=e^{-i n x}$.

Compare to Distributional Derivative, Definition 11.10.

Example 7.60. Weak Derivative of $f(x)=|x|$
Notes 1/26/11

$$
f(x)=|x| \quad-\pi<x<\pi
$$

$f \in C(\mathbb{T})$, but its standard derivative $f^{\prime} \notin C(\mathbb{T})$ because $f^{\prime}(0)$ and $f^{\prime}(\pi)$ don't exist. We shall see if $g=f^{\prime}$ (weak derivative) exists. We want:

$$
\begin{aligned}
\int g \phi d x & =-\int f \phi^{\prime} d x \\
& =-\int_{0}^{\pi} x \phi^{\prime} d x+\int_{-\pi}^{0} x \phi^{\prime} d x \\
& =-\left.x \phi\right|_{0} ^{\pi}+\int_{0}^{\pi} \phi d x+\left.x \phi\right|_{-\pi} ^{0}-\int_{-\pi}^{0} \phi d x \\
& ==\pi \phi(\pi)+\underline{\pi} \phi(-\pi)+\int_{-\pi}^{\pi} \operatorname{sgn} x \phi d x
\end{aligned}
$$

We conclude that $\int f \phi x=-\int g \phi d x \forall \phi \in C^{\infty}(\mathbb{T})$ if $g(x)=\operatorname{sgn} x$.

Example 7.61. Weak Derivative of $f(x)=\operatorname{sgn} x$
Notes $1 / 26 / 11$

$$
\begin{aligned}
\int h \phi d x & =-\int g \phi^{\prime} d x \\
& =-\int_{0}^{\pi} \phi^{\prime} d x+\int_{-\pi}^{0} \phi^{\prime} d x \\
& =-[\phi(\pi)-\phi(0)]+[\phi(0)-\phi(-\pi)] \\
& =2[\phi(0)-\phi(\pi)]
\end{aligned}
$$

There is no such $h \in L^{1}$. To see this, take $\phi=\frac{1}{2 \pi} e^{-i n x} \in C^{\infty}(\mathbb{T})$.

$$
\hat{h}(n)=\frac{1}{\pi}\left[1-e^{i n \pi}\right]=\left\{\begin{array}{cl}
\frac{2}{\pi} & n \text { odd } \\
0 & n \text { even }
\end{array}\right.
$$

This contradicts the Riemann-Lebesge Lemma, and therefore there is no such $h \in L^{1}$.

## Proposition 7.62.

Notes $1 / 26 / 11$
$f$ is weakly differentiable with $f \in L^{1}$ iff it is absolutely continuous.

## Definition 7.63. Absolutely Continuous

http://en.wikipedia.org/wiki/Absolute_continuity\#Absolute_continuity_of_functions
$f$ is absolutely continuous if it has a derivative $f^{\prime}$ a.e., the derivative is Lebesgue integrable, and

$$
f(x)=f(a)+\int_{a}^{x} f^{\prime}(t) d t
$$

## Theorem 7.64.

Notes 1/26/11

If $f$ is weakly differentiable with weak derivative $g=f^{\prime} \in L^{1}(\mathbb{T})$, then

$$
\hat{g}(n)=i n \hat{f}(n)
$$

## $\underline{\text { Proof }}$

$$
\hat{g}(n)=\frac{1}{2 \pi} \int g(x) e^{-i n x} d x=-\frac{1}{2 \pi} \int f(x) e^{-i n x} d x=i n \hat{f}(n)
$$

## Proposition 7.65.

Notes $1 / 26 / 11$

A function $f \in L^{2}(\mathbb{T})$ has a weak derivative $g \in L^{2}(\mathbb{T})$ iff

$$
\sum_{n \in \mathbb{Z}} n^{2}|\hat{f}(n)|^{2}<\infty
$$

and then

$$
g(x)=\sum_{n \in \mathbb{Z}} i n \hat{f}(n) e^{i n x}
$$

Definition 7.66. Sobolev Space (2)
Notes $1 / 26 / 11$

The Sobolev space $W^{1, p}(\mathbb{T}), 1 \leq p \leq \infty$, consists of all functions $f: \mathbb{T} \rightarrow \mathbb{C}$ s.t. $f \in L^{p}(\mathbb{T})$, $f^{\prime} \in L^{p}(\mathbb{T})$. If $p=2$, we write $W^{1,2}(\mathbb{T})=H^{1}(\mathbb{T})$ (where the H is because it is a Hilbert space).

A function $f \in H^{1}(\mathbb{T})$ iff

$$
\sum_{n \in \mathbb{Z}}\left(1+n^{2}\right)|\hat{f}(n)|^{2}<\infty
$$

and

$$
\begin{aligned}
\|f\|_{H^{1}} & =\left(\int|f|^{2} d x+\int\left|f^{\prime}\right|^{2} d x\right)^{1 / 2} \\
& =\left(\|f\|_{L^{2}}^{2}+\left\|f^{\prime}\right\|_{L^{2}}^{2}\right)^{1 / 2} \\
& =\left(2 \pi \sum_{n \in \mathbb{Z}}\left(1+n^{2}\right)|\hat{f}(n)|^{2}\right)^{1 / 2}
\end{aligned}
$$

## Theorem 7.67. Sobolev Embedding Theorem

Notes 1/26/11

If $f \in H^{1}(\mathbb{T})$ then $f \in C(\mathbb{T})$ and

$$
\|f\|_{\infty} \leq C\|f\|_{H^{1}}
$$

$J: H^{1} \rightarrow C$ (Embedding), $f \mapsto f$.

## $\underline{\text { Proof }}$

$$
\begin{aligned}
\sum_{n \in \mathbb{Z}}|\hat{f}(n)| & =\sum_{n \in \mathbb{Z}} \frac{1}{\left(1+n^{2}\right)^{1 / 2}}\left(1+n^{2}\right)^{1 / 2}|\hat{f}(n)| \\
& \leq\left(\sum_{n \in \mathbb{Z}} \frac{1}{\left(1+n^{2}\right)^{1 / 2}}\right)\left(\sum_{n \in \mathbb{Z}}\left(1+n^{2}\right)|\hat{f}(n)|\right) \\
& \leq C\|f\|_{H^{1}}
\end{aligned}
$$

It follows that $f \in C(\mathbb{T})$ because the Fourier series converges uniformly to $f$ (see Theorem 7.26) and

$$
\|f\|_{\infty} \leq \sum_{n \in \mathbb{Z}}|\hat{f}(n)| \leq C\|f\|_{H^{1}}
$$

### 7.7 Chapter Summary

This chapter explores the spaces $L^{p}(\mathbb{T}), p \in[1, \infty)$, with special attention given to the Hilbert space $L^{2}(\mathbb{T})$. These spaces are the completion of $C(\mathbb{T})$ with respect to the $L^{p}$-norm; thus, $C(\mathbb{T})$ is dense in $L^{p}(\mathbb{T})$ for $p \in[1, \infty)$. Since $\mathbb{T}$ has finite Lebesgue measure, we can use Hölder's Inequality to show that for $p>q,\|\cdot\|_{p} \geq\|\cdot\|_{q}$, which implies that $L^{p}(\mathbb{T}) \subset L^{q}(\mathbb{T})$. We define the convolution of two functions and what it means for a family of functions to be an approximate identity, and we use these tools to prove the Weierstrass Approximation Theorem, which says that the trigonometric polynomials are dense in $C(\mathbb{T})$ with respect to the uniform norm. Since uniform convergence implies $L^{2}$ convergence, it follows that the functions $e_{n}(x)=\frac{1}{\sqrt{2 \pi}} e^{i n x}$ form an orthonormal basis for $L^{2}(\mathbb{T})$. Thus, for all $f \in L^{2}(\mathbb{T})$, we have that

$$
f(x)=\sum_{n=-\infty}^{\infty} \hat{f}_{n} e^{i n x}
$$

where the equality is in the $L^{2}$ sense. A result from Carleson tells us that the Fourier series of $f$ converges pointwise to $f$ a.e.

Next we explore some properties of Fourier series and Fourier coefficients. Let $f, g \in L^{2}(\mathbb{T})$. We use the density of $C(\mathbb{T})$ in $L^{2}(\mathbb{T})$ to prove the Convolution Theorem, which allows us to express the Fourier coefficients of $f * g$ in terms of those of $f$ and $g:(\widehat{(f * g})_{n}=\sqrt{2 \pi} \hat{f}_{n} \hat{g}_{n}$. Parseval's Theorem allows us to compute $\langle f, g\rangle$ using the Fourier coefficients of $f$ and $g:\langle f, g\rangle=\sum_{n=-\infty}^{\infty} \overline{\hat{f}_{n}} \hat{g}_{n}$.

Now we examine the Fourier series of differentiable functions. Using integration by parts, we show that

$$
\hat{f}_{n}^{\prime}=i n \hat{f}_{n} .
$$

This gives us the concept of a weak derivative, since the derivative of $f$ may not be continuous; e.g. $f(x)=|x|$. We define the Sobolev space $H^{k}(\mathbb{T})$ as the space of $L^{2}(\mathbb{T})$ functions with $k$ weak derivatives. And since the boundary terms on $\mathbb{T}$ vanish, we have that $\left\langle f^{\prime}, g\right\rangle=-\left\langle f, g^{\prime}\right\rangle$ for $f, g \in H^{1}(\mathbb{T})$. Thus, we may define the weak derivative of a function using integration by parts: $g \in L^{1}(\mathbb{T})$ is the weak derivative of $f \in L^{1}(\mathbb{T})$ if

$$
\int_{\mathbb{T}} f \phi^{\prime} d x=-\int_{\mathbb{T}} g \phi d x \quad \forall \phi \in C^{\infty}(\mathbb{T}) .
$$

Finally, we prove a special case of the Sobolev Embedding Theorem: if $f \in H^{k}(\mathbb{T})$ for $k>1 / 2$, then $f \in C(\mathbb{T})$.
In addition, Hunter briefly discussed $L^{1}(\mathbb{T})$. We can define the Fourier series of an $L^{1}$ function, but we cannot guarantee that it converges to the function. Our main result is the Riemann-Lebesgue Lemma, which says that the Fourier coefficients of an $L^{1}$ function decay to zero as $n \rightarrow \infty$. Hunter then discussed 3 kernels: the Dirichlet kernel (standard summation), Fejér kernel (Cesáro summation), and Poisson kernel (Abel summation). These kernels are related to the concept of approximate identities, and we convolve the kernels with a function $f$. He covered harmonic functions, and our main result is that we can use the Poisson kernel to solve the two-dimensional Laplace equation.

## 11 Distributions and the Fourier Transform

### 11.1 Periodic Distributions

## Definition 11.1. Test Functions

Notes $1 / 28 / 11$ and http://en.wikipedia.org/wiki/Distribution_\(mathematics\) and Hunter's Notes page 51

We define our space of test functions as:
$\mathcal{D}(\mathbb{T})=C^{\infty}(\mathbb{T})$ with the following topology:
$\varphi_{n} \rightarrow \varphi \in \mathcal{D}$ if $\varphi_{n}^{(k)} \rightarrow \varphi^{(k)}$ uniformly for all $k=0,1,2, \ldots$ Note that this topology is not obtained from any norm, but rather it is derived.

## Definition 11.2. Distribution

Notes $1 / 28 / 11$ and Hunter's Notes page 51

A distribution is a continuous linear functional, $T$, that maps a set of test functions, $\mathcal{D}(\mathbb{T})$, onto the set of complex numbers. The space of distributions is denoted by $\mathcal{D}^{\prime}(\mathbb{T})$. For $T \in \mathcal{D}^{\prime}(\mathbb{T}), \varphi \in \mathcal{D}(\mathbb{T})$, we write:

$$
\langle T, \varphi\rangle=T(\varphi)
$$

$\mathcal{D}^{\prime}(\mathbb{T})$ is the topological dual space of the distributions on $\mathbb{T}$ (i.e. $\mathcal{D}(\mathbb{T})$ ), with the topology defined as follows: $T_{n} \rightharpoonup T$ in $\mathcal{D}^{\prime}$ if $\left\langle T_{n}, \varphi\right\rangle \rightarrow\langle T, \varphi\rangle$ in $\mathbb{C} \forall \varphi \in \mathcal{D}$.

$$
\begin{aligned}
& T: \mathcal{D}(\mathbb{T}) \rightarrow \mathbb{C} \\
& \text { Linear: }\langle T, \lambda \varphi+\mu \psi\rangle=\lambda\langle T, \varphi\rangle+\mu\langle T, \psi\rangle \\
& \text { Continuous: If } \varphi_{n} \rightarrow \varphi \in \mathcal{D}, \text { then }\left\langle T, \varphi_{n}\right\rangle \rightarrow\langle T, \varphi\rangle \in \mathbb{C}
\end{aligned}
$$

Compare Distributional Convergence, $T_{n} \rightharpoonup T$ in $\mathcal{D}^{\prime}$ if $\left\langle T_{n}, \varphi\right\rangle \rightarrow\langle T, \varphi\rangle$, to Weak Convergence (Definition 8.41): $x_{n} \rightharpoonup x$ if $\left\langle x_{n}, y\right\rangle \rightarrow\langle x, y\rangle \quad \forall y \in \mathcal{H}$.

## Definition 11.3. Seminorm

Notes 1/28/11

Our topology on $\mathcal{D}$ is obtained from a countable family of seminorms:

$$
\|\varphi\|_{k}=\sup _{x \in \mathbb{T}}\left|\varphi^{(k)}(x)\right|, \quad k=0,1,2, \ldots
$$

A seminorm has the same properties as a norm except that it may assign length zero to nonzero vectors.

## Example 11.4. Seminorms

Notes 1/28/11

$$
d(\varphi, \psi)=\sum_{k=0}^{\infty} \frac{1}{2^{k}} \frac{\|\varphi-\psi\|_{k}}{1+\|\varphi-\psi\|_{k}}
$$

- This is not a norm because you can't pull out a constant
- This turns $\mathcal{D}$ into a Fréchet space (a complete, metrizable topological vector space topology defined by a countable family of seminorms)
- We could instead use norms to define the topology on $\mathcal{D}(\mathbb{T})$ :

$$
\|\varphi\|_{C^{k}}=\sum_{j=0}^{k}\|\varphi\|_{j}
$$

## Remark 11.5.

Notes 1/28/11

Note that the differentiation operator

$$
D: \mathcal{D}(\mathbb{T}) \rightarrow \mathcal{D}(\mathbb{T}), \quad D(\varphi)=\varphi^{\prime}
$$

is continuous: if $\varphi_{n} \rightarrow \varphi \in \mathcal{D}$, then $D \varphi_{n} \rightarrow D \varphi \in \mathcal{D}$. This is because there are inifinitely many semi-norms.

## Example 11.6. Regular Distribution

page 292 and Notes 1/28/11

If $f: \mathbb{T} \rightarrow \mathbb{C}$ is integrable, $f \in L^{1}(\mathbb{T})$, define

$$
\begin{aligned}
& T_{f}: \mathcal{D}(\mathbb{T}) \rightarrow \mathbb{C} \\
& T_{f}(\varphi)=\int_{\mathbb{T}} f \varphi d x
\end{aligned}
$$

$\left|T_{f}(\varphi)\right| \leq \sup |\varphi| \cdot \int|f| d x<\infty$, so $T_{f}$ is well-defined. It is a distribution because it satisfies:

1. Linearity: (1) $T_{f}(\varphi+\psi)=\int f(\varphi+\psi) d x=T_{f}(\varphi)+T_{f}(\psi)$. (2) $T_{f}(c \varphi)=c T_{f}(\varphi)$
2. Continuity: If $\varphi_{n} \rightarrow 0$ in $\mathcal{D}$, then $\left|T_{f}\left(\varphi_{n}\right)\right| \leq \sup \left|\varphi_{n}\right|\|f\|_{L^{1}} \rightarrow 0$ as $n \rightarrow \infty$. So $T_{f}\left(\varphi_{n}\right) \rightarrow 0$ and $T_{f}$ is continuous.

We identify $f$ with $T_{f}$. Thus, $L^{1}(\mathbb{T}) \subset D^{1}(\mathbb{T})$.
We call $T_{f}$ a regular distribution. A regular distribution is a distribution that is given by the integration of a test function $\varphi$ against a function $f$.

## Definition 11.7. Principal Value Distribution

page 293

A principal value distribution is a singular distribution, denoted by p.v. $(1 / x)$, and its action on a test function $\varphi$ is given by

$$
\text { p.v. } \frac{1}{x}(\varphi)=\lim _{\epsilon \rightarrow 0^{+}} \int_{|x|>\epsilon} \frac{\varphi(x)}{x} d x
$$

## Example 11.8.

Notes $1 / 28 / 11$

Consider the periodic $\delta$-function (actually a distribution, not a function).

$$
\begin{aligned}
& \langle\delta, \varphi\rangle=\varphi(0) \\
& \langle\delta, \varphi+\psi\rangle=(\varphi+\psi)(0)=\varphi(0)+\psi(0)=\langle\delta, \varphi\rangle+\langle\delta, \psi\rangle \\
& \langle\delta, c \varphi\rangle=c\langle\delta, \varphi\rangle
\end{aligned}
$$

$\varphi_{n} \rightarrow 0$ implies $\varphi_{n}(0) \rightarrow 0$, and therefore $\delta$ is a continuous linear functional.
$\delta$ is not regular. Proof:

- Suppose $\langle\delta, \varphi\rangle=\int f \varphi d x$ for some $f \in L^{1}$.
- Consider $\varphi_{n}(x)=\left[\frac{1+\cos x}{2}\right]^{n}$
- $\left\langle\delta, \varphi_{n}\right\rangle=1 \forall n$, but $\int f \varphi_{n} d x \rightarrow 0$ as $n \rightarrow \infty$ by the Lebesge-Dominated Convergence Theorem if $f \in L^{1}$
- Thus, there is no function $f \in L^{1}$ such that $\int f \varphi d x=\varphi(0)$


## Example 11.9.

Notes 1/28/11

$$
\text { Let } T_{n}=\left\{\begin{aligned}
\frac{1}{2} n & |x| \leq \frac{1}{n} \\
0 & \frac{1}{n} \leq|x| \leq \pi
\end{aligned}\right.
$$

Then $\int_{-\pi}^{\pi} T_{n} d x=1 \forall n$. Claim: $\left\langle T_{n}, \varphi\right\rangle=\frac{n}{2} \int_{1 / n}^{1 / n} \varphi(x) \rightarrow \varphi(0)$ as $n \rightarrow \infty$. Proof:

$$
\begin{aligned}
\left|\frac{n}{2} \int_{-1 / n}^{1 / n} \varphi(x) d x-\varphi(0)\right| & =\frac{n}{2}\left|\int_{-1 / n}^{1 / n}[\varphi(x)-\varphi(0)] d x\right| \\
& \leq \frac{n}{2}\left[\sup _{|x| \leq 1 / n}|\varphi(x)-\varphi(0)|\right] \cdot \frac{2}{n} \\
& \leq \sup _{|x| \leq 1 / n}|\varphi(x)-\varphi(0)| \rightarrow 0 \text { as } n \rightarrow \infty
\end{aligned}
$$

## Definition 11.10. Distributional Derivative

Every distribution $T \in \mathcal{D}^{\prime}(\mathbb{T})$ has a distributional derivative $T^{\prime} \in \mathcal{D}(\mathbb{T})$ that is given by

$$
\left\langle T^{\prime}, \phi\right\rangle=-\left\langle T, \phi^{\prime}\right\rangle \quad \forall \phi \in \mathcal{D}(\mathbb{T})
$$

Compare to Weak Derivative (2), Definition 7.59.

## Definition 11.11. Motivation for Distributional Derivatives

Notes 1/31/11

Suppose $f \in C^{\infty}$ is a smooth function. Consider $T_{f^{\prime}}$ :

$$
\left\langle T_{f^{\prime}}, \varphi\right\rangle=\int f^{\prime} \varphi d x=-\int f \varphi^{\prime} d x=-\left\langle T_{f}, \varphi^{\prime}\right\rangle
$$

Want: $\left(T_{f^{\prime}}\right)=\left(T_{f}\right)^{\prime}$

This defines the distributional derivative.

1. Linearity: $\left\langle T^{\prime}, a \varphi+b \psi\right\rangle=-\left\langle T,(a \varphi+b \psi)^{\prime}\right\rangle=-\left\langle T, a \varphi^{\prime}+b \psi^{\prime}\right\rangle=-a\left\langle T, \varphi^{\prime}\right\rangle-b\left\langle T, \psi^{\prime}\right\rangle=$ $a\left\langle T^{\prime}, \varphi\right\rangle+b\left\langle T^{\prime}, \psi\right\rangle$
2. Continuity: Suppose $\varphi_{n} \rightarrow \varphi$ in $\mathcal{D}$. Consider $\left\langle T^{\prime}, \varphi\right\rangle$.
$\left\langle T^{\prime}, \varphi_{n}\right\rangle=-\left\langle T, \varphi_{n}^{\prime}\right\rangle \rightarrow-\left\langle T, \varphi^{\prime}\right\rangle=\left\langle T^{\prime}, \varphi\right\rangle$, because $T$ is continuous on $\mathcal{D}$ and $D: \varphi \rightarrow \varphi^{\prime}$ is continuous on $\mathcal{D}$

Example 11.12.
Notes $1 / 31 / 11$

$$
\begin{aligned}
f(x) & =|x|, \quad|x| \leq \pi \\
f^{\prime}(x) & =\operatorname{sgn} x=g(x)
\end{aligned}
$$

Compute the distributional derivative of $g$ :

$$
\begin{aligned}
\left\langle g^{\prime}, \varphi\right\rangle & =-\left\langle g, \varphi^{\prime}\right\rangle \\
& =-\int_{0}^{\pi} \varphi^{\prime} d x+\int_{-\pi}^{0} \varphi^{\prime} d x \\
& =-[\varphi(\pi)-\varphi(0)]+[\varphi(0)-\varphi(\pi)] \\
& =2 \varphi(0)-2 \varphi(\pi) \\
& =2\left\langle\delta_{0}, \varphi\right\rangle-2\left\langle\delta_{\pi}, \varphi\right\rangle \\
& =\left\langle 2 \delta_{0}-2 \delta_{\pi}, \varphi\right\rangle \\
g^{\prime} & =2 \delta_{0}-2 \delta_{\pi} \\
& =2\left(\delta-\tau_{\pi} \delta\right)
\end{aligned}
$$

Where $\tau_{\pi}$ means translation by $\pi$ and $\delta_{a}$ is the $\delta$-"function" supported at $a$ :

$$
\left\langle\delta_{a}, \varphi\right\rangle=\varphi(a)
$$

Example 11.13.
Notes $1 / 31 / 11$

Compute $\delta^{\prime}$ :

$$
\left\langle\delta^{\prime}, \varphi\right\rangle=-\left\langle\delta, \varphi^{\prime}\right\rangle=-\varphi^{\prime}(0)
$$

## Definition 11.14. Fourier Coefficients

Notes $1 / 31 / 11$

If $T \in \mathcal{D}^{\prime}(\mathbb{T})$, define $\hat{T}(n)=\frac{1}{2 \pi}\left\langle T, e^{-i n x}\right\rangle$.

## Example 11.15.

Notes $1 / 31 / 11$

Compute the Fourier coefficients of $\delta$ :

$$
\begin{aligned}
& \hat{\delta}(n)=\frac{1}{2 \pi}\left\langle\delta, e^{-i n x}\right\rangle=\frac{1}{2 \pi} e^{0}=\frac{1}{2 \pi} \\
& \delta(x)=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} e^{i n x}
\end{aligned}
$$

## Remark 11.16.

$1 / 31 / 11$

There are 3 contexts in which to look at Fourier series:

- Continuous functions $\Rightarrow$ converge uniformly
- $L^{2}$ functions $\Rightarrow$ converge in $L^{2}$
- Distribution functions $\Rightarrow$ converge in the distributional sense

Example 11.17.
Notes $1 / 31 / 11$

$$
P_{r}(x)=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{i n x}
$$

Formally, as $r \rightarrow 1^{-}, P_{r}(x) \rightharpoonup \frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{i n x}=\delta(x)$

## Theorem 11.18.

Notes $1 / 31 / 11$
$\varphi \in \mathcal{D}$ iff $(\hat{\varphi}(n))$ is rapidly decreasing, i.e.

$$
|n|^{k} \hat{\varphi}(n) \rightarrow 0 \text { as } n \rightarrow \infty \forall k \geq 0
$$

and the Fourier series of $\varphi$ converges to $\varphi$ in $\mathcal{D}$.

Compare to Corollary 7.55.

## Proof

- $\varphi \in C^{k} \Rightarrow|n|^{k} \hat{\varphi}(n) \rightarrow 0$ by the Riemann-Lebesgue Lemma, so if $\varphi \in C^{\infty}$, then the $\hat{\varphi}(n)$ are rapidly decreasing
- Sobolev Embedding Theorem: If $\hat{\varphi}(n)$ is rapidly decreasing, then $\varphi \in H^{k}(\mathbb{T}) \forall k$ implies that

$$
\sum\left(1+n^{2}\right)|\hat{\varphi}(n)|^{2}<\infty
$$

- Hence, $\varphi \in C^{k-1}(\mathbb{T}) \forall k$. So $\varphi \in C^{\infty}$.
- Similarly, $\sum_{|n| \leq N} \hat{\varphi}(n) e^{i n x} \rightarrow \varphi$ in $H^{k} \forall k$
- So $\sum_{|n| \leq N} \hat{\varphi}(n) e^{i n x} \rightarrow \varphi$ in $C^{k-1} \forall k$
- So $\sum_{|n| \leq N} \hat{\varphi}(n) e^{i n x}$ converges in $\mathcal{D}$


## Definition 11.19. $S(\mathbb{Z})$

Notes 2/2/11
$S(\mathbb{Z})$ is the space of rapidly decreasing sequences, $\left(c_{n}\right)$, such that

$$
\lim _{n \rightarrow \infty}|n|^{k} c_{n}=0 \quad \forall k=0,1,2, \ldots
$$

## Remark 11.20.

Notes 2/2/11
$\mathcal{F}: C^{\infty}(\mathbb{T}) \rightarrow S(\mathbb{Z})$
$\mathcal{F}: \varphi \rightarrow(\hat{\varphi}(n))$
If $\varphi \in C^{\infty}(\mathbb{T})$, then $S_{N} \varphi=\sum_{|n| \leq N} \hat{\varphi}(n) e^{i n x} \rightarrow \varphi$ in $\mathcal{D}$.
If $T \in \mathcal{D}^{\prime}(\mathbb{T})$, then $\hat{T}(n)=\frac{1}{2 \pi}\left\langle T, e^{-i n x}\right\rangle$

## Proposition 11.21.

Notes 2/2/11

$$
\widehat{T}^{\prime}(n)=i n \hat{T}(n)
$$

See Proposition 7.52 and Definition 7.56.

Proof.

$$
\begin{aligned}
\widehat{T}^{\prime}(n) & =\frac{1}{2 \pi}\left\langle T^{\prime}, e^{-i n x}\right\rangle=-\frac{1}{2 \pi}\left\langle T,\left(e^{-i n x}\right)^{\prime}\right\rangle=i n \cdot \frac{1}{2 \pi}\left\langle T, e^{-i n x}\right\rangle \\
& =\operatorname{in} \hat{T}(n)
\end{aligned}
$$

## Definition 11.22. Slow Growth

Notes 2/2/11

A sequence $\left(c_{n}\right)$ has slow growth if there exist $k, M$ such that $\left|c_{n}\right| \leq M\left(1+n^{2 k}\right)^{1 / 2} \forall n$.
Equivalently, $\left|c_{n}\right| \leq M|n|^{k} \forall n \neq 0$.

Lemma 11.23.
Notes 2/2/11

If $T \in \mathcal{D}^{\prime}$, then $(\hat{T}(n))$ has slow growth.

Proof. If $T \in \mathcal{D}^{\prime}$ then $T$ has some finite order $k$ such that

$$
|\langle T, \varphi\rangle| \leq C\|\varphi\|_{C^{k}}
$$

Then

$$
|\hat{T}(n)|=\left|\left\langle T, e^{-i n x}\right\rangle\right| \leq C\left\|e^{-i n x}\right\|_{C^{k}} \leq C\left(1+n^{2 k}\right)^{1 / 2}
$$

Example 11.24. Weierstrass Nowhwere Differentiable Function
Notes 2/2/11

$$
f(x)=\sum_{n=1}^{\infty} \frac{1}{2^{n}} \cos \left(3^{n} x\right)
$$

$\sum \frac{1}{2^{n}}<\infty$, so $f \in \mathcal{A}(\mathbb{T})$.

$$
f^{\prime}(x)=\sum_{n=1}^{\infty} \frac{3^{n}}{2^{n}} \sin \left(3^{n} x\right)
$$

$f$ is nowhere differentiable, although it does have a distributional derivative.

## Theorem 11.25.

Notes 2/2/11

If $T \in \mathcal{D}^{\prime}(\mathbb{T})$ and $S_{N} T=\sum_{|n| \leq N} \hat{T}(n) e^{i n x} \in C^{\infty}(\mathbb{T})$, then $S_{N} T \rightharpoonup T$ in $\mathcal{D}^{\prime}$ as $N \rightarrow \infty$.
Ex: $\delta(x)=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} e^{i n x}$

Proof.

$$
\begin{aligned}
\left\langle S_{N} T, \varphi\right\rangle & =\left\langle\sum_{|n| \leq N} \hat{T}(n) e^{i n x}, \varphi\right\rangle=\sum_{|n| \leq N}\left\langle\hat{T}(n) e^{-i n x}, \varphi\right\rangle=\sum_{|n| \leq N} \hat{T}(n) \int e^{i n x} \varphi(x) d x \\
& =2 \pi \sum_{|n| \leq N} \hat{T}(n) \hat{\varphi}(-n)=2 \pi \sum_{|n| \leq N}\left\langle T, e^{-i n x}\right\rangle \cdot \frac{1}{2 \pi} \hat{\varphi}(-n)=\left\langle T, \sum_{|n| \leq N} \hat{\varphi}(-n) e^{-i n x}\right\rangle \\
& =\left\langle T, S_{N} \varphi\right\rangle \rightarrow\langle T, \varphi\rangle \text { as } n \rightarrow \infty
\end{aligned}
$$

So $S_{N} T \rightarrow T$ as $N \rightarrow \infty$.

## Theorem 11.26.

Notes 2/2/11

If $\left(c_{n}\right)$ is a sequence of slow growth, $\left(c_{n}\right) \in S^{\prime}(\mathbb{Z})$, then there exists a distribution $T$ such that $\hat{T}(n)=c_{n}$.

Proof. Define $T$ by

$$
\langle T, \varphi\rangle=2 \pi \sum_{n \in \mathbb{Z}} c_{n} \hat{\varphi}(-n)
$$

Remark 11.27.
Notes $2 / 2 / 11$
$\mathcal{F}: f \mapsto \hat{f}(n)$

$$
\begin{array}{rlr}
\mathcal{D}(\mathbb{T})=C^{\infty}(\mathbb{T}) & \leftrightarrow S(\mathbb{Z}) \\
& C(\mathbb{T}) \supset \mathcal{A}(\mathbb{T}) \leftrightarrow \ell^{\prime}(\mathbb{Z}) \\
L^{2}(\mathbb{T}) & \leftrightarrow \ell^{2}(\mathbb{Z}) \\
& L^{1}(\mathbb{T}) \rightarrow C_{0}(\mathbb{Z}) \\
\mathcal{D}^{\prime}(\mathbb{T}) & \leftrightarrow S^{\prime}(\mathbb{Z}) &
\end{array}
$$

- $C^{\infty} \subset L^{2}(\mathbb{T}) \subset \mathcal{D}^{\prime}(\mathbb{T})$
- $S(\mathbb{Z}) \subset \ell^{2}(\mathbb{Z}) \subset S^{\prime}(\mathbb{Z})$


## 8 Bounded Linear Operators on a Hilbert Space

### 8.1 Orthogonal Projections

## Definition 8.1. Direct Sum

page 187 and Notes $2 / 4 / 11$

If $M$ and $N$ are subspaces of a linear space $X$ such that every $x \in X$ can be written uniquely as $x=y+z$ with $y \in M$ and $z \in N$, then we say that $X=M \oplus N$ is the direct sum of $M$ and $N$, and we call $N$ a complementary subspace of $M$ in $X$. The decomposition $x=y+z$ is unique if and only if $M \cap N=\{0\}$.

Definition 8.2. Projection, Idempotent, Self-Adjoint
page $187 \& 188$ and Notes $2 / 4 / 11$

Given a direct sum decomposition, $X=M \oplus N$, define the projection $P: X \rightarrow X$ onto $M$ along $N$ by

$$
P(m+n)=m, \quad m \in M, \quad n \in N
$$

All projections are linear and idempotent, meaning that $P^{2}=P$, because

$$
P^{2}(m+n)=P(m)=m
$$

## Theorem 8.3.

page 188 and Notes 2/4/11

Any linear map $P: X \rightarrow X$ with $P^{2}=P$ is a projection. Specifically, it is the projection onto ran $P$ along ker $P$.

Proof.

- $x=P(x)+(x-P(x))$
- $P^{2}(x)=P(x) \quad \Rightarrow \quad P(x) \in \operatorname{ran} P$
- $P(x-P(x))=P x-P^{2} x=P x-P x=0 \quad \Rightarrow \quad x-P(x) \in \operatorname{ker} P$
- Suppose $x \in \operatorname{ker} P \cap \operatorname{ran} P$
$-x \in \operatorname{ran} P \quad \Rightarrow \quad x=P y$
$-x \in \operatorname{ker} P \quad \Rightarrow \quad 0=P x=P^{2} y=P y=x=0$
- Thus, $x=0$, and ker $P \cap \operatorname{ran} P=\{0\}$
- Thus, $X=\operatorname{ran} P \oplus \operatorname{ker} P$


## Remark 8.4. Bounded Projections

Notes 2/4/11

Question: Given a projection $P: X \rightarrow X, X$ a Banach space, when can we say that $P$ is bounded?

Answer: We need ran $P$ closed and complemented by a closed subspace $N=\operatorname{ker} P$

Note: The kernel of a bounded operator is always closed; the range need not be.

## Definition 8.5. Orthogonal Projections, Self-Adjoint

Notes $2 / 4 / 11$ and $2 / 7 / 11$

Let $\mathcal{H}$ be a Hilbert space and let $M \subset \mathcal{H}$ be a closed linear subspace. Then by the Projection Theorem,

$$
\mathcal{H}=M \oplus M^{\perp}, \quad M^{\perp}=\{y \in \mathcal{H} \mid y \perp m \forall m \in M\}
$$

We define the orthogonal projection $P: \mathcal{H} \rightarrow \mathcal{H}$ onto $M$ along $M^{\perp}$.

An orthogonal projection $P$ on a Hilbert space $\mathcal{H}$ is

- Idempotent: $P^{2}=P$
- Self-Adjoint: $\langle x, P y\rangle=\langle P x, y\rangle$

Proof. To see that a projection $P$ on a Hilbert space $\mathcal{H}$ is self-adjoint, let

$$
x=m+n, \quad y=p+q, \quad \text { where } \quad m, p \in M, \quad n, q \in N
$$

Compute:

$$
\begin{aligned}
& \langle x, P y\rangle=\langle m+n, p\rangle=\langle m, p\rangle+\langle n, p\rangle=\langle m, p\rangle \\
& \langle P x, y\rangle=\langle m, p+q\rangle=\langle m, p\rangle+\langle m, q\rangle=\langle m, p\rangle
\end{aligned}
$$

Lemma 8.6.
page 188 and Notes 2/7/11

If $P$ is a nonzero othogonal projection then $\|P\|=1$

Proof.

$$
\|P x\|^{2}=\langle P x, P x\rangle=\left\langle x, P^{2} x\right\rangle=\langle x, P x\rangle \leq\|x\|\|P x\|
$$

Either $\|P x\|=0$ or $\|P x\| \leq\|x\|$. Since $\|P x\| \neq 0 \forall x$, it must be the case that $\|P x\| \leq\|x\|$. Then

$$
\|P\|=\sup \frac{\|P x\|}{\|x\|} \leq 1
$$

If $P \neq 0$, then there exists $y \in \mathcal{H}$ such that $P y \neq 0$. Setting $x=P y$ in the previous equation yields

$$
\|P\| \geq \frac{\|P \cdot P x\|}{\|P x\|}=1
$$

So $\|P\|=1$.

## Theorem 8.7.

page 189 and Notes 2/7/11

If $P$ is an orthogonal projection, then $\mathcal{H}=M \oplus M^{\perp}=\operatorname{ran} P \oplus \operatorname{ker} P$, where $M=\operatorname{ran} P$ and $M^{\perp}=\operatorname{ker} P$ are closed subspaces. Conversely, if $M$ is any closed subspace of $\mathcal{H}$, then there exists an orthogonal projection with $M=\operatorname{ran} P$ and $M^{\perp}=\operatorname{ker} P$.

## Example 8.8. Even $\mathcal{G}$ Odd Functions

page 189 and Notes 2/7/11

Let $\mathcal{H}=L^{2}(\mathbb{R})$ and let

$$
\begin{aligned}
& M=\text { space of even functions, } f(-x)=f(x) \\
& N=\text { space of odd functions, } f(-x)=-f(x)
\end{aligned}
$$

$M \perp N$, since $\int \bar{f} g d x=0$ for $f$ odd, $g$ even. Define

- Even Projection: $P: \mathcal{H} \rightarrow \mathcal{H}$ onto $M, P f(x)=\frac{1}{2}[f(x)+f(-x)]$
- Odd Projection: $Q: \mathcal{H} \rightarrow \mathcal{H}$ onto $N, Q f(x)=\frac{1}{2}[f(x)-f(-x)]$
- Note: $Q=I-P$

Check that $P$ is self-adjoint:

$$
\langle P f, g\rangle=\int_{\mathbb{R}} \frac{1}{2} \overline{[f(x)+f(-x)]} g(x) d x=\int_{\mathbb{R}} \frac{1}{2} \bar{f}(x) g(x)+\frac{1}{2} \bar{f}(x) g(-x) d x=\langle f, P g\rangle
$$

## Example 8.9.

Notes 2/7/11

Let $\mathcal{H}=L^{2}(\mathbb{T})$. Define $P f=\frac{1}{2 \pi} \int_{\mathbb{T}} f d x, P: \mathcal{H} \rightarrow \mathcal{H}$.

$$
\text { Given: } f=\sum_{n \in \mathbb{Z}} \hat{f}(n) e^{i n x}
$$

Then: $P f=\hat{f}(0)$

- Idempotent: $P^{2}=P$ since $P f$ is a constant, and $P 1=1$
- Self-Adjoint: $\langle P f, g\rangle=\int \overline{\left[\frac{1}{2 \pi} \int f d x\right]} g d x=\frac{1}{2 \pi} \int \bar{f} d x \int g d x=\langle f, 1\rangle \cdot \frac{1}{2 \pi} \int g d x=\langle f, P g\rangle$
ran $P=$ constant functions $=<1>($ space spanned by 1$)$
$\operatorname{ker} P=$ functions with zero mean (i.e. $\hat{f}(0)=0$ )
ran $P \perp$ ker $P$


## Example 8.10. Fourier Projections

Notes 2/7/11

We can define the orthogonal projection of $f$ onto the $N$ th partial sum of its Fourier series:

$$
P_{N} f=\sum_{|n| \leq N} \hat{f}(n) e^{i n x}
$$

Similarly, we can define the projection onto the positive $n$ part of its Fourier series:

$$
\begin{gathered}
P f=\sum_{n=0}^{\infty} \hat{f}(n) e^{i n x} \\
(I-P) f=\sum_{n=-\infty}^{-1} \hat{f}(n) e^{i n x}
\end{gathered}
$$

Example 8.11.
page 189 and Notes 2/7/11

Let $\mathcal{H}=L^{2}(\mathbb{R})$. If $A \subset \mathbb{R}$ is some Lebesgue measurable set, define

$$
\chi_{A}(x)= \begin{cases}1 & x \in A \\ 0 & x \notin A\end{cases}
$$

Then

$$
P_{A} f=\chi_{A} f
$$

is an orthogonal projection of $L^{2}(\mathbb{R})$ onto the subspace of functions with support contained in $\bar{A}$.

### 8.2 The Dual of a Hilbert Space

Theorem 8.12. Riesz Representation Theorem
page 191 and Notes 2/7/11

Given: a Hilbert space $\mathcal{H}$, its dual space $\mathcal{H}^{*}=\mathcal{B}(\mathcal{H}, \mathbb{C})$ (the set of bounded linear maps $\varphi: \mathcal{H} \rightarrow \mathbb{C}$ with $\left.\|\varphi\|_{\mathcal{H}^{*}}=\sup \frac{|\varphi(x)|}{\|x\|}<\infty\right)$.

Every $\varphi \in \mathcal{H}^{*}$ can be given by $\varphi(x)=\langle y, x\rangle$ for some $y \in \mathcal{H}$, and $\|\varphi\|=\|y\|$. Conversely, every $y \in \mathcal{H}$ corresponds to a $\varphi \in \mathcal{H}^{*}$. The map $J: \varphi \mapsto y$ is an isometric, antilinear isomorphism of $\mathcal{H}^{*}$ onto $\mathcal{H}$.

$$
\begin{array}{ll}
\text { Antilinear: } & J(\varphi+\psi)=J(\varphi)+J(\psi) \\
& J(\lambda \varphi)=\bar{\lambda} J(\varphi)
\end{array}
$$

Proof.

- Suppose $\varphi \in \mathcal{H}^{*}$. We want to find $y \in \mathcal{H}$ such that $\varphi(x)=\langle y, x\rangle$
- Suppose $\varphi \neq 0$. Then $\operatorname{ker} \varphi \neq \mathcal{H}$ and $\operatorname{ker} \varphi$ is closed because $\varphi$ is bounded
- There exists $z \in(\operatorname{ker} \varphi)^{\perp}$ (by the Projection Theorem)
- Consider $P: \mathcal{H} \rightarrow \mathcal{H}, P x=\frac{\varphi(x)}{\varphi(z)} P z$. Claim: this is an orthogonal projection.
- Idempotent: $P^{2} x=P\left(\frac{\varphi(x)}{\varphi(z)} z\right)=\frac{\varphi(x)}{\varphi(z)} P z=\frac{\varphi(x)}{\varphi(z)} z \quad$ (since $\left.P z=z\right)$
- Self-Adjoint: $\langle x, P y\rangle=\langle P x, y\rangle$
- H$=\operatorname{ran} P \oplus \operatorname{ker} P, \quad$ ran $P=\langle z\rangle, \quad \operatorname{ker} P=\operatorname{ker} \varphi$
- $x \in \mathcal{H}, \quad x=\alpha z+w, \quad w \in \operatorname{ker} \varphi, \quad \alpha=\frac{\langle z, x\rangle}{\|z\|^{2}}$
- $\varphi(x)=\alpha \varphi(z)=\frac{\langle z, x\rangle}{\|z\|^{2}} \varphi(z)=\langle y, x\rangle, \quad y=\frac{\bar{\varphi}(z)}{\|z\|^{2}} z$


### 8.3 The Adjoint of an Operator

## Definition 8.13. Adjoint

page 193 and Notes 2/9/11

Given a bounded linear map $A \in \mathcal{B}(\mathcal{H})$, its adjoint $A^{*} \in \mathcal{B}(\mathcal{H})(\leftarrow$ proved in Proposition 8.15) is the linear map that satisfies

$$
\langle x, A y\rangle=\left\langle A^{*} x, y\right\rangle \quad \forall x, y \in \mathcal{H}
$$

Remark 8.14. Adjoint: Existence and Uniqueness
page 193 and Notes 2/9/11

To define $A^{*}$ such that $A^{*} x=z$, consider $\varphi_{x}: \mathcal{H} \rightarrow \mathbb{C}, \varphi_{x}(y)=\langle x, A y\rangle$. Then

$$
\begin{aligned}
\left\|\varphi_{x}(y)\right\| & \leq\|x\|\|A y\| \leq\|x\|\|A\|\|y\| \\
\left\|\varphi_{x}\right\| & \leq\|A\|\|x\|
\end{aligned}
$$

So $\varphi_{x}$ is a bounded linear functional. By the Riesz Representation Theorem, there is a unique $z \in \mathcal{H}$ such that

$$
\varphi_{x}(y)=\langle z, y\rangle
$$

Define $A^{*} x=z$. Then

$$
\begin{aligned}
& \langle x, A y\rangle=\varphi_{x}(y)=\langle z, y\rangle=\left\langle A^{*} x, y\right\rangle \\
& \langle x, A y\rangle=\left\langle A^{*} x, y\right\rangle \quad \forall x, y \in \mathcal{H}
\end{aligned}
$$

## Proposition 8.15.

Notes 2/9/11

If $A \in \mathcal{B}(\mathcal{H})$ then $A^{*} \in \mathcal{B}(\mathcal{H})$ and
(1) $\quad\left\|A^{*}\right\|=\|A\|$
(2) $\|A\|^{2}=\left\|A^{*} A\right\|$
(See also Corollary 8.34.)

Proof.

$$
\begin{aligned}
& \left\|A^{*}\right\|=\sup _{\|x\|=1}\left\|A^{*} x\right\| \quad \text { (See Lemma } 8.26 \text { in the book) } \\
& \quad=\sup _{\|x\|=\|y\|=1}\left|\left\langle y, A^{*} x\right\rangle\right|=\sup _{\|x\|=\|y\|=1}|\langle A y, x\rangle|=\sup _{\|y\|=1}\|A y\|=\|A\| \\
& \|A\|^{2}=\sup _{\|x\|=1}\|A x\|^{2}=\sup _{\|x\|=1}|\langle A x, A x\rangle|=\sup _{\|x\|=1}\left|\left\langle x, A^{*} A x\right\rangle\right| \\
& \quad \leq\left\|A^{*} A\right\| \quad \text { (See Corollary } 8.27 \text { in the book) } \\
& \left\|A^{*} A\right\| \leq\left\|A^{*}\right\|\|A\|=\|A\|^{2} \\
& \left\|A^{*} A\right\|
\end{aligned}=\|A\|^{2} \quad l
$$

Remark 8.16.
Notes 2/9/11
$\mathcal{B}(\mathcal{H})$ is a $C^{*}$-algebra.

$$
\|A B\| \leq\|A\|\|B\| \quad *: \mathcal{B}(\mathcal{H}) \rightarrow \mathcal{B}(\mathcal{H}), * *=\text { identity } \quad\left\|A^{*}\right\|=\|A\|
$$

## Remark 8.17. Generalizations

Notes 2/9/11

1. Given: $A: \mathcal{H} \rightarrow K, A^{*}: K \rightarrow \mathcal{H}$, where $\mathcal{H}, K$ are Hilbert spaces.
$\langle x, A y\rangle_{K}=\left\langle A^{*} x, y\right\rangle_{H} \quad \forall y \in \mathcal{H}, x \in K$
$A^{*}$ is the Hilbert space adjoint.
2. Given: $A: X \rightarrow Y, A^{\prime}: Y^{\prime} \rightarrow X^{\prime}$, where $X, Y$ are Banach spaces and $X^{\prime}$ is the dual space of $X$.
$\langle\psi, A x\rangle_{Y \times Y^{\prime}}=\left\langle A^{\prime} \psi, x\right\rangle_{X \times X^{\prime}} \forall x \in X, \psi \in Y^{\prime}$
$A^{\prime}$ is the dual operator or Banach space adjoint.

Example 8.18.
page 193 and Notes $2 / 9 / 11$ and Notes $2 / 11 / 11$

Let $\mathcal{H}=\mathbb{C}^{n}$. Then $A: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ is given by a matrix $\left(a_{i j}\right)$.

$$
\begin{aligned}
& y_{i}=\sum_{j=1}^{n} a_{i j} x_{j}, \quad x=\left(x_{1}, \ldots, x_{n}\right), \quad y=\left(y_{1}, \ldots, y_{n}\right) \\
&\langle x, y\rangle=\sum_{i=1}^{n} x_{i} y_{i} \\
&\langle x, A y\rangle=\sum_{i=1}^{n} \overline{x_{i}}\left(\sum_{j=1}^{n} a_{i j} y_{j}\right)=\sum_{j=1}^{n}\left(\overline{\sum_{i=1}^{n} \overline{a_{i j}} x_{i}}\right) y_{j} \\
&=\left\langle A^{*} x, y\right\rangle \\
& \text { If } z=A^{*} x \\
& z_{j}=\sum_{i=1}^{n} \overline{a_{i j}} x_{i}=\sum_{j=1}^{n} \overline{a_{j i}} x_{j}
\end{aligned}
$$

- $A^{*}$ has matrix $\left(\overline{a_{j i}}\right)$, which is the conjugate transpose of $\left(a_{i j}\right)$
- $\left(A^{*} A\right)$ is Hermitian, positive definite
- $\left(A^{*} A\right)^{*}=\left(A^{*} A\right)^{*}=A^{*} A$
- $\left\langle x, A^{*} A x\right\rangle=\langle A x, A x\rangle \geq 0$
- $A^{*} A$ has orthogonal eigenvectors that form a basis of $\mathbb{C}^{n}$ with eigenvalues $\mu_{1}, \mu_{2}, \ldots, \mu_{n} \geq 0$
- $\left\|A^{*} A\right\|=\max _{1 \leq j \leq n}\left|\mu_{j}\right|=\sigma\left(A^{*} A\right)=$ the spectral radius of $A^{*} A$
- $\|A\|=\sqrt{\sigma\left(A^{*} A\right)}$

Example 8.19.
page 194 and Notes 2/9/11

Let $\mathcal{H}=L^{2}([0,1]),\langle f, g\rangle=\int_{0}^{1} \overline{f(x)} g(x) d x$.
Define the integral operator $K: L^{2}([0,1]) \rightarrow L^{2}([0,1])$ by

$$
K f(x)=\int_{0}^{1} k(x, y) f(y) d y, \quad k:[0,1] \times[0,1] \rightarrow \mathbb{C}
$$

(Note: $k(x, y)$ is the kernel of the integral operator $K$. It is not related to the null space.) Ex: Assume that $k$ is Hilbert-Schmidt: $k$ is measurable on $[0,1] \times[0,1]$ and

$$
\begin{aligned}
\|K\|^{2} & \leq \int_{0}^{1} \int_{0}^{1}|k(x, y)|^{2} d x d y<\infty \\
\langle f, K g\rangle & =\int_{0}^{1} \overline{f(x)}\left(\int_{0}^{1} k(x, y) g(y) d y\right) d x \\
& =\int_{0}^{1}\left(\overline{\int_{0}^{1} f(x) \overline{k(x, y)} d x}\right) g(y) d y \\
& =\left\langle K^{*} f, g\right\rangle
\end{aligned}
$$

Since

$$
\begin{aligned}
K^{*} f(y) & =\int_{0}^{1} \overline{k(x, y)} f(y) d x \\
K^{*} f(x) & =\int_{0}^{1} \overline{k(y, x)} f(y) d y
\end{aligned}
$$

Thus, $K^{*}$ is an integral operator with conjugate transpose level of $k$.

Example 8.20.
page 194 and Notes 2/9/11

Recall the right and left shift operators, respectively:

$$
S\left(x_{1}, x_{2}, x_{3}, \ldots\right)=\left(0, x_{1}, x_{2}, x_{3}, \ldots\right) \quad T\left(x_{1}, x_{2}, x_{3} \ldots\right)=\left(x_{2}, x_{3}, \ldots\right)
$$

$T$ is the adjoint of $S$, i.e. $T=S^{*}$. Also, $S=T^{*}$.

## Example 8.21. Solvability of Linear Equations

Notes $2 / 11 / 11$

Consider $A: \mathcal{H} \rightarrow \mathcal{H}, A x=y$. Suppose for some $y \in \mathcal{H}$ we have a solution for $x \in \mathcal{H}$.

Let $z \in \operatorname{ker} A^{*}$. Then

$$
\langle z, A x\rangle=\left\langle A^{*} z, x\right\rangle=\langle z, y\rangle
$$

Thus, a necessary condition for solvability is that $y \perp z \forall z \in \operatorname{ker} A^{*}$, i.e. $y \perp \operatorname{ker} A^{*}$.

## Theorem 8.22.

page 194 and Notes $2 / 11 / 11$

If $A \in \mathcal{B}(\mathcal{H})$, then $\mathcal{H}=\overline{\operatorname{ran} A} \oplus\left(\operatorname{ker} A^{*}\right)$, and

$$
\overline{\operatorname{ran} A}=\left(\operatorname{ker} A^{*}\right)^{\perp} \quad \operatorname{ker} A=\left(\operatorname{ran} A^{*}\right)^{\perp}
$$

Proof. From Example 8.21, if $y \in \operatorname{ran} A$ then $y \in\left(\operatorname{ker} A^{*}\right)^{\perp}$.

$$
\begin{aligned}
& \operatorname{ran} A \subset\left(\operatorname{ker} A^{*}\right)^{\perp} \\
& \overline{\operatorname{ran} A} \subset\left(\operatorname{ker} A^{*}\right)^{\perp} \quad \text { since orthogonal complements are closed }
\end{aligned}
$$

If $y \in(\operatorname{ran} A)^{\perp}$ then

$$
\begin{aligned}
\langle A x, z\rangle & =0 \forall x \in \mathcal{H} \\
\left\langle x, A^{*} y\right\rangle & =0 \forall x \in \mathcal{H}
\end{aligned}
$$

This implies that $A^{*} y=0$, so $y \in \operatorname{ker} A^{*}$.

$$
\begin{gathered}
(\operatorname{ran} A)^{\perp} \subset \operatorname{ker} A^{*} \\
\overline{\operatorname{ran} A}=(\operatorname{ran} A)^{\perp \perp} \supset\left(\operatorname{ker} A^{*}\right)^{\perp}
\end{gathered}
$$

Corollary 8.23.
page 195 and Notes $2 / 11 / 11$

If $A \in \mathcal{B}(\mathcal{H})$ has closed range (ran $A$ is a closed linear subspace), then $A x=y$ is solvable iff $y \perp \operatorname{ker} A^{*}$.

## Example 8.24.

Notes 2/11/11

If $\mathcal{H}$ is finite dimensional, or $A$ has finite rank, then $\operatorname{ran} A$ is closed and Corollary 8.23 applies.

## Example 8.25.

page 196 and Notes $2 / 11 / 11$

Recall the left $(T)$ and right $(S)$ shift operators. $S^{*}=T, T^{*}=S$.

1. $\mathcal{H}=\overline{\operatorname{ran} S} \oplus \operatorname{ker} S^{*}=\overline{\operatorname{ran} S} \oplus \operatorname{ker} T$
2. $\mathcal{H}=\overline{\operatorname{ran} T} \oplus \operatorname{ker} T^{*}=\overline{\operatorname{ran} T} \oplus \operatorname{ker} S$

- $\operatorname{ran} S=\left\{\left(x_{1}, x_{2}, \ldots\right) \in \ell^{2} \mid x_{1}=0\right\} \quad$ - $\operatorname{ran} T=\ell^{2}(\mathbb{N})$
- $\operatorname{ker} S=\{0\}$
- $\operatorname{ker} T=\left\{\left(x_{1}, 0,0,0, \ldots\right) \mid x_{1} \in \mathbb{C}\right\}$
$S x=y$ is solvable iff $y \perp \operatorname{ker} T$, and the solution is unique.
$T x=y$ is solvable for all $y \in \ell^{2}(\mathbb{N})$, but the solution is not unique.


## Example 8.26.

Notes 2/11/11

Let $\mathcal{H}=\ell^{2}(\mathbb{N}), A\left(x_{1}, x_{2}, \ldots, x_{n}, \ldots\right)=\left(x, \frac{1}{2} x_{2}, \ldots, \frac{1}{n} x_{n}, \ldots\right)$.

$$
[A]=\left(\begin{array}{cccccc}
1 & & & & & \\
& \frac{1}{2} & & & & \\
& & \frac{1}{3} & & & \\
& & & \ddots & & \\
& & & & \frac{1}{n} & \\
& & & & & \ddots
\end{array}\right), \quad A^{*}=A \text { (self-adjoint) }
$$

$\operatorname{ker} A=\operatorname{ker} A^{*}=\{0\}$
$\mathcal{H}=\overline{\operatorname{ran} A} \oplus \operatorname{ker} A$

Given $y=\left(y_{1}, y_{2}, \ldots\right) \in \ell^{2}(\mathbb{N})$, does there exist $x=\left(x_{1}, x_{2}, \ldots\right) \in \ell^{2}(\mathbb{N})$ such that $A x=y ?$
$x \in \ell^{2}(\mathbb{N}) \Leftrightarrow \sum n^{2}\left|y_{n}\right|<\infty$
$\operatorname{ran} A=\left\{\left.\left(x_{1}, x_{2}, \ldots\right) \in \ell^{2}(\mathbb{N})\left|\sum n^{2}\right| x_{n}\right|^{2}<\infty\right\}$
$\operatorname{ran} A \neq \mathcal{H}$, so $A$ is not onto.
$\mathrm{Ex}: M=\left\{\left(x_{1}, x_{2}, \ldots, x_{N}, 0,0, \ldots\right)\right\} \subset \operatorname{ran} A$
$M$ is dense in $\ell^{2}(\mathbb{N})$, so $\overline{\operatorname{ran} A}=\ell^{2}(\mathbb{N}), \quad \ell^{2}(\mathbb{N})=\overline{\operatorname{ran} A} \oplus \operatorname{ker} A^{*}$

Consider: $A x=y, y \perp \operatorname{ker} A^{*}=\operatorname{ker} A=\{0\} \forall y \in \ell^{2}(\mathbb{N})$. This is not solvable for every $y \in \ell^{2}(\mathbb{N})$, only for $y \in \operatorname{ran} A$, and $\operatorname{ran} A$ is a dense, non-closed subspace of $\ell^{2}(\mathbb{N})$.

### 8.4 Self-Adjoint and Unitary Operators

## Definition 8.27. Self-Adjoint

page 197 and 2/14/11

A bounded operator $A: \mathcal{H} \rightarrow \mathcal{H}$ on a Hilbert space $\mathcal{H}$ is self-adjoint if $A^{*}=A$.

Equivalently, $A$ is self-adjoint iff

$$
\langle x, A y\rangle=\langle A x, y\rangle \quad \forall x, y \in \mathcal{H}
$$

## Example 8.28. Self-Adjoint Operators

Notes 2/14/11

1. $A: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n},[A]^{*}=[A]$
$A: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n},[A]^{T}=[A]$
2. $\mathcal{H}=L^{2}(\mathbb{R})$. Suppose $a: \mathbb{R} \rightarrow \mathbb{C}$ is bounded and measurable. Define $M: \mathcal{H} \rightarrow \mathcal{H}, M f=a f$. $\|M f\|_{2} \leq\|a\|_{\infty}\|f\|_{2}$.
$M^{*} f=\bar{a} f, M^{*}=M$ if $a: \mathbb{R} \rightarrow \mathbb{R}$.
3. Orthogonal projections: $P^{2}=P=P^{*}$ (self-adjoint)
4. Given $T \in \mathcal{B}(\mathcal{H}), A=T^{*} T$ is self-adjoint.
$T=A+i B, A=\frac{1}{2}\left(T^{*}+T\right), B=\frac{1}{2 i}\left(T^{*}-T\right)$
$A^{*}=A, B^{*}=B$
5. The shift operators are NOT self-adjoint because $S^{*}=T \neq S$

## Definition 8.29. Bilinear Forms, Sesquilinear

page 197 and Notes 2/14/11

Let $A: \mathcal{H} \rightarrow \mathcal{H}$ be a bounded linear operator. We define the bilinear form $a: \mathcal{H} \times \mathcal{H} \rightarrow \mathbb{C}$ by

$$
a(x, y)=\langle x, A y\rangle
$$

We say that $a$ is sesquilinear because

$$
\begin{aligned}
& a(x, \lambda y+\mu z)=\lambda a(x, y)+\mu a(x, z) \\
& a(\lambda x+\mu y, z)=\bar{\lambda} a(x, z)+\bar{\mu} a(x, z)
\end{aligned}
$$

Definition 8.30. Hermitian Symmetric 63 Symmetric
page 197 and Notes 2/14/11

Suppose $A$ is self-adjoint. Then

$$
\begin{gathered}
\langle x, A y\rangle=\langle A x, y\rangle=\overline{\langle y, A x\rangle} \\
a(x, y)=\overline{a(x, y)}
\end{gathered}
$$

We say that $a$ is Hermitian symmetric. In the real case, we have $a(x, y)=a(y, x)$, and we say that this is symmetric.

## Definition 8.31. Quadratic Form

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Given $A: \mathcal{H} \rightarrow \mathcal{H}$, we define the quadratic form $q: \mathcal{H} \rightarrow \mathbb{C}$ by

$$
q(x)=\langle x, A x\rangle=a(x, x)
$$

If $A$ is self-adjoint, then $a(x, x)=\overline{a(x, x)}$, so $a(x, x)$ is real for all $x \in \mathcal{H}$.

Definition 8.32. Positive, Positive Definite
page 198 and Notes 2/14/11

A self-adjoint operator $A$ is positive or positive definite if $\langle x, A x\rangle=a(x, x)>0$ for all $x \in \mathcal{H}, x \neq 0$.

## Theorem 8.33.

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If $A$ is self-adjoint then

$$
\|A\|=\sup _{x \neq 0} \frac{|\langle x, A x\rangle|}{\|x\|^{2}}=\sup _{\|x\|=1}|\langle x, A x\rangle|
$$

Note: compare this to $\|A\|=\sup _{\|x\|=1}|\langle A x, A x\rangle|^{1 / 2}$ (see part 2 of Proposition 8.15).

Proof.

$$
|\langle x, A x\rangle| \leq\|x\|\|A x\| \leq\|A\|\|x\|^{2} \quad \text { (Cauchy-Schwarz) }
$$

Let $\alpha=\sup _{\|x\| \neq 0} \frac{|\langle x, A x\rangle|}{\|x\|^{2}} \leq\|A\|$. Then $|\langle x, A x\rangle| \leq \alpha\|x\|^{2} \leq\|A\|\|x\|^{2}$. The parallelogram law states that

$$
\langle x, A y\rangle=\frac{1}{4}\{\langle x+y, A(x+y)\rangle-\langle x-y, A(x-y)\rangle-i\langle x+i y, A(x+i y)\rangle+i\langle x-i y, A(x-i y)\rangle\}
$$

In general,

$$
\|A\|=\sup _{\|x\|=\|y\|=1}|\langle x, A y\rangle|
$$

and this does not require self-adjoint. If $A$ is self-adjoint, the first 2 terms in the parallelogram law expression are real and the last 2 are imaginary. We can multiply $y$ by $e^{i \theta}$ so that $e^{i \theta}\langle x, A y\rangle=\langle x, A z\rangle$ is real, where $z=y e^{i \theta}$. Then we have

$$
\begin{aligned}
e^{i \theta}\langle x, A y\rangle & =\langle x, A z\rangle \\
& =\frac{1}{4}\{\langle x+z, A(x+z)\rangle-\langle x-z, A(x-z)\rangle\} \\
|\langle x, A y\rangle| & \leq \frac{1}{4}|\langle x+z, A(x+z)\rangle|+\frac{1}{4}|\langle x-z, A(x-z)\rangle| \\
& \leq \frac{\alpha}{4}\left(\|x+z\|^{2}+\|x-z\|^{2}\right) \\
& \leq \frac{\alpha}{2}\left(\|x\|^{2}+\|z\|\right) \quad(\text { by the parallelogram rule (not law)) } \\
\|A\| & \leq \sup _{\|x\|=\|y\|=1}|\langle x, A y\rangle| \leq \frac{\alpha}{2}\left(\|x\|^{2}+\|y\|^{2}\right) \leq \frac{\alpha}{2}(1+1)=\alpha
\end{aligned}
$$

Corollary 8.34.
page 199

If $A$ is a bounded operator on a Hilbert space then $\left\|A^{*} A\right\|=\|A\|^{2}$. If $A$ is self-adjoint, then $\left\|A^{2}\right\|=\|A\|^{2}$.

The proof follows directly from Proposition 8.15.

Definition 8.35. Unitary Operators
pages $199 \& 200$ and Notes $2 / 14 / 11$

An operator $U: \mathcal{H} \rightarrow \mathcal{H}$ is unitary if

$$
U^{*} U=U U^{*}=I, \quad \text { i.e. } U^{*}=U^{-1}
$$

Note that

$$
\langle U x, U y\rangle=\left\langle U^{*} U x, y\right\rangle=\langle x, y\rangle
$$

so $U$ preserves norms and inner products. Furthermore, if $\left\{e_{n} \mid n \in \mathbb{N}\right\}$ is an orthonormal basis of $\mathcal{H}$, then so is $\left\{U e_{n} \mid n \in \mathbb{N}\right\}$.

## Example 8.36.

Notes 2/14/11

1. $U: \mathbb{C}^{2} \rightarrow \mathbb{C}^{2}$ with matrix

$$
[U]=\left(\begin{array}{rc}
a & b \\
-\bar{b} & \bar{a}
\end{array}\right), \quad|a|^{2}+\left|b^{2}\right|=1, \quad a, b \in \mathbb{C}
$$

In the real case, $a=\cos \theta, b=\sin \theta$, and $U$ is rotation by $\theta$.
2. The right shift operator $S$ on $\ell^{2}(\mathbb{N})$ is not unitary because

$$
S^{*}=T, \quad S^{*} S=I, \quad S S^{*}=P \neq I
$$

3. If $A^{*}=A$ then $U=e^{i A}$ is unitary, where

$$
\begin{aligned}
e^{i A} & =I+(i A)+\cdots+\frac{1}{n!}(i A)^{n}+\ldots \\
U^{*} & =e^{-i A} \\
U^{*} U & =I
\end{aligned}
$$

## Example 8.37. Quantum Mechanics

Notes 2/14/11

In quantum mechanics we have the Hamiltonian operator $H$, with $H^{*}=H$. We also have $U(t)=$ $e^{i t H}, U: \mathcal{H} \rightarrow K, U^{*}: K \rightarrow \mathcal{H} . U$ is unitary if $U^{*} U=I_{H}$ and $U U^{*}=I_{K}$. We say that 2 Hilbert spaces are isometric if they are unitarily equivalent.

Example 8.38.
page 201 and Notes $2 / 14 / 11$

$$
\begin{aligned}
& \mathcal{F}: L^{2}(\mathbb{T}) \rightarrow \ell^{2}(\mathbb{Z}) \quad \text { is unitary } \\
& \mathcal{F} f=\hat{f}, \quad \hat{f}(n)=\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{T}} f(x) e^{-i n x} d x
\end{aligned}
$$

## Definition 8.39. Normal Operators

Notes 2/16/11

If $T: \mathcal{H} \rightarrow \mathcal{H}$ is a bounded linear operator on a Hilbert space $\mathcal{H}$, then $T$ is normal if

$$
\left[T^{*}, T\right] \equiv T^{*} T-T T^{*}=0 \quad \text { i.e. } \quad T^{*} T=T T^{*}
$$

Self-adjoint and unitary operators are normal.

## Example 8.40.

Notes 2/16/11

1. Self-adjoint and unitary operators are normal
2. The shift operators on $\ell^{2}(\mathbb{N})$ are not normal
3. Any multiplication operator is normal

$$
\begin{aligned}
& M: L^{2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R}) \\
& (M f)(x)=m(x) f(x), \quad m \in L^{\infty}(\mathbb{R}) \\
& M^{*} f=\bar{m} f \\
& M^{*} M f=\bar{m} m f=m \bar{m} f=M M^{*} f
\end{aligned}
$$

Special cases
(a) If $m$ is real-valued then $M=M^{*}$, so $M$ is self-adjoint. For
(b) For $M$ to be unitary, we must have $m=e^{i \theta}$.

### 8.6 Weak Convergence in a Hilbert Space

## Definition 8.41. Weak Convergence

page 204 and Notes 2/16/11

A sequence $\left(x_{n}\right)$ in a Hilbert space $\mathcal{H}$ converges weakly to $x \in \mathcal{H}$, written $x_{n} \rightharpoonup x$, if

$$
\left\langle x_{n}, y\right\rangle \rightarrow\langle x, y\rangle \quad \forall y \in \mathcal{H}
$$

Compare to Distributional Convergence (Definition 11.2): $T_{n} \rightharpoonup T$ in $\mathcal{D}^{\prime}$ if $\left\langle T_{n}, \varphi\right\rangle \rightarrow\langle T, \varphi\rangle$.

## Definition 8.42. Strong Convergence

Notes 2/16/11

We write strong (norm) convergence as $x_{n} \rightarrow x$ if $\left\|x_{n}-x\right\| \rightarrow 0$.

## Remark 8.43. Weak vs. Strong Convergence

Notes 2/16/11

If $x_{n} \rightarrow x$, then $x_{n} \rightharpoonup x$ because

$$
\left|\left\langle x_{n}, y\right\rangle-\langle x, y\rangle\right| \leq\left\|x_{n}-x\right\|\|y\| \quad \text { (Cauchy-Schwarz) }
$$

In a finite dimensional space, the converse is true, but this is not the case in infinite dimensional spaces.

Weak convergence $=$ component-wise convergence

## Example 8.44.

page 204 and Notes 2/16/11

Let $\mathcal{H}$ be a separable Hilbert space and let $\left\{e_{n} \mid n \in \mathbb{N}\right\}$ be a separable orthonormal basis. Then $e_{n} \rightharpoonup 0$ as $n \rightarrow \infty$ because

$$
\left\langle e_{n}, y\right\rangle=y_{n} \rightarrow 0 \quad \text { as } \quad n \rightarrow \infty \quad \text { because } \quad \sum\left|y_{n}\right|^{2}<\infty
$$

But $\left(e_{n}\right)$ doesn't converge strongly because

$$
\left\|e_{n}-e_{m}\right\|=\sqrt{2} \quad \forall n \neq m
$$

and so the sequence is not Cauchy and hence not convergent.

## Example 8.45.

Notes 2/16/11

Define an unbounded sequence $\left(x_{n}\right)$ by $x_{n}=n e_{n}$. We know that

$$
\left\langle x_{n}, e_{m}\right\rangle \rightarrow 0 \quad \Rightarrow \quad\left\langle x_{n}, y\right\rangle \rightarrow 0 \quad \text { as } \quad n \rightarrow \infty \quad \forall y=\sum_{m=1}^{m} c_{m} e_{m}
$$

Let $y_{1}=\sum \frac{1}{m} e_{m}$. Then

$$
\left\langle x_{n}, y\right\rangle=\frac{1}{n} \cdot n=1 \quad \forall n
$$

Let $y_{2}=\sum \frac{1}{m^{3 / 4}} e_{m} \in \mathcal{H}$. Then

$$
\left\langle x_{n}, y\right\rangle=\frac{1}{n^{3 / 4}} \cdot n \rightarrow 0
$$

Thus, $\left(x_{n}\right)$ does not converge weakly.

Theorem 8.46. Uniform Boundedness Theorem
page 204

Suppose that $\left\{\varphi_{n}: X \rightarrow \mathbb{C} \mid n \in \mathbb{N}\right\}$ is a set of functionals on a Banach space $X$ such that the set of complex numbers $\left\{\varphi_{n}(x) \mid n \in \mathbb{N}\right\}$ is bounded for each $x \in X$. Then $\left\{\left\|\varphi_{n}\right\| \mid n \in \mathbb{N}\right\}$ is bounded.

## Theorem 8.47.

Notes 2/16/11

If $x_{n} \rightharpoonup x$ then $\left\{\left\|x_{n}\right\| \mid n \in \mathbb{N}\right\}$ is bounded.

Proof. Define $\varphi_{n}: \mathcal{H} \rightarrow \mathbb{C}$ by $\varphi_{n}(y)=\left\langle x_{n}, y\right\rangle$. Then $\varphi_{n} \in \mathcal{H}^{*}$. By the uniform boundedness theorem (Theorem 8.46),

$$
\begin{aligned}
& \left|\varphi_{n}(y)\right| \leq M \quad \forall y \in \mathcal{H}, n \in \mathbb{N} \\
& \left\{\left|\varphi_{n}(y)\right| \mid n \in \mathbb{N}\right\} \text { is bounded for each } y \in \mathcal{H}, \text { so }\left\{\left\|\varphi_{n}\right\| \mid n \in \mathbb{N}\right\} \text { is bounded }
\end{aligned}
$$

## Theorem 8.48.

page 205 and Notes $2 / 16 / 11$

Let $D \subset \mathcal{H}$ be a dense subset. Then $x_{n} \rightharpoonup x$ iff
(a) $\left\{\left\|x_{n}\right\| \mid n \in \mathbb{N}\right\}$ is bounded
(b) $\left\langle x_{n}, y\right\rangle \rightarrow\langle x, y\rangle \quad \forall y \in D$

## Proposition 8.49.

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$$
\text { If } x_{n} \rightharpoonup x \text {, then }\|x\| \leq \liminf _{n \rightarrow \infty}\left\|x_{n}\right\|
$$

Proof.

$$
\begin{gathered}
\|x\|^{2}=\langle x, x\rangle=\lim _{n \rightarrow \infty}\left\langle x_{n}, x\right\rangle \leq\|x\| \liminf _{n \rightarrow \infty}\left\|x_{n}\right\| \\
\left\langle x_{n}, x\right\rangle \leq\left\|x_{n}\right\|\|x\| \quad \text { (Cauchy-Schwarz) }
\end{gathered}
$$

Note: if $a_{n} \leq b_{n}, a_{n} \rightarrow a$, then $a \leq \liminf b_{n}$.

$$
\left\|x_{n}-x\right\|^{2}=\left\langle x_{n}-x, x_{n}-x\right\rangle=\left\|x_{n}\right\|^{2}-\left\langle x, x_{n}\right\rangle-\left\langle x_{n}, x\right\rangle+\|x\|^{2}
$$

If $x_{n} \rightharpoonup x$, then $\left\|x_{n}\right\| \rightarrow\|x\|$, and

$$
\left\|x_{n}-x\right\|^{2} \rightarrow\|x\|^{2}-\langle x, x\rangle-\langle x, x\rangle+\|x\|^{2}=0
$$

Example 8.50. Example for Proposition 8.49
Notes 2/16/11

$$
\begin{array}{lll}
x_{1}=e_{1} & x_{n} \rightharpoonup 0 \\
x_{2}=2 e_{2} \\
x_{3}=e_{3} & & \\
x_{4}=2 e_{4} & \liminf _{n \rightarrow \infty}=1
\end{array} \quad\left\|x_{n}\right\|= \begin{cases}1 & n \text { odd } \\
2 & n \text { even }\end{cases}
$$

...

## Example 8.51. Weak Convergence $\nRightarrow$ Strong Convergence

Notes 2/16/11

## (a) Oscillation:

(1) Let $\mathcal{H}=L^{2}(\mathbb{T}), f_{n}(x)=e^{i n x} \rightharpoonup 0$ as $n \rightarrow \infty$

Proof. $\left\|f_{n}\right\|=\sqrt{2 \pi}$ is bounded, and $\left\langle e^{i n x}, \varphi\right\rangle \rightarrow 0$ as $n \rightarrow \infty$ for all trig polynomials $\varphi$, and the trig polynomials are dense in $L^{2}(\mathbb{T})$.
(2) Let $\mathcal{H}=L^{2}(\mathbb{R})$. Recall that $C_{C}^{\infty}(\mathbb{R}) \subset L^{2}(\mathbb{R})$ are the smooth functions with compact support, and they are dense in $L^{2}(\mathbb{R})$. Then $f_{n} \rightharpoonup f$ iff
i. $\|f\| \leq M$ (bounded)
ii. $\int f_{n} \varphi d x \rightarrow \int f \varphi d x \forall \varphi \in C_{C}^{\infty}(\mathbb{R})$

Consider $f_{n}(x)=\psi(x) \sin (n \pi x)$, where $\psi \in C_{C}^{\infty}(\mathbb{R}) \cap L^{2}(\mathbb{R})$. Then $f_{n} \rightharpoonup 0$ as $n \rightarrow \infty$, but $f_{n} \nrightarrow 0$ as $n \rightarrow \infty$. (See proof below)
(b) Concentration: Consider

$$
f_{n}(x)=\left\{\begin{array}{rl}
n^{1 / 2} & 0<x<\frac{1}{n} \\
0 & \text { otherwise }
\end{array}\right.
$$

i. $\left\|f_{n}\right\|^{2}=\int_{0}^{1 / 2}\left(n^{1 / 2}\right)^{2} d x=1$
ii. $\forall \varphi \in C_{C}^{\infty}(\mathbb{R}), \quad\left|\int f_{n} \varphi d x\right|=\left|n^{1 / 2} \int_{0}^{1 / n} \varphi d x\right| \leq n^{1 / 2} \cdot \frac{1}{n}\|\varphi\|_{\infty} \rightarrow 0$ as $n \rightarrow \infty$ So $f_{n} \rightharpoonup 0$ as $n \rightarrow \infty$
Does $f_{n}$ converge strongly to 0 ? No, because $\left\|f_{n}\right\|=1 \forall n$. (See below for more details)
(c) Escape to Infinity:

$$
f_{n}(x)= \begin{cases}1 & n<x<n+1 \\ 0 & \text { otherwise }\end{cases}
$$

i. $\left\|f_{n}\right\|_{L^{2}}=1$, so $f_{n}$ is bounded.
ii. $\int f_{n} \varphi d x \rightarrow 0$ as $n \rightarrow \infty \forall \varphi \in C_{C}^{\infty}(\mathbb{R})$

Thus, $f_{n} \rightharpoonup 0$, but $f_{n} \nrightarrow 0$ because $\left\|f_{n}\right\|=1 \forall n$.

Proof. (a2)
i. $\left\|f_{n}\right\|^{2}=\int \psi^{2}(x) \sin ^{2}(n \pi x) d x \leq \int \psi^{2}(x) d x \leq\|\psi\|^{2}$
ii. Suppose $\varphi \in C_{C}(\mathbb{R})$.

$$
\begin{aligned}
\int f_{n}(x) \varphi(x) d x & =\int \psi(x) \sin (n \pi x) \varphi(x) d x \\
& \left.=\int \frac{\cos (n \pi x)}{n \pi}[\varphi(x) \psi(x)]^{\prime} d x \quad \quad \text { (IBP, no boundary terms because } \varphi \in C_{C}(\mathbb{R})\right) \\
\left|\int f_{n} \varphi d x\right| & \leq \frac{1}{n \pi} \int(|\varphi \psi|)^{\prime} d x \\
& \leq \frac{c}{n}
\end{aligned}
$$

So $\int f_{n} \varphi d x \rightarrow 0$ as $n \rightarrow \infty$, and thus $f_{n} \rightharpoonup f$.

Does $\left(f_{n}\right)$ converge strongly? i.e., does $f_{n} \rightarrow 0$ ? (see Remark 8.52)
If $\psi \neq 0$, then

$$
\left\|f_{n}\right\|^{2}=\int \psi^{2}(x) \sin ^{2}(n \pi x) d x=\int \psi^{2}(x) \cdot \frac{1}{2}[1-\cos (2 n \pi x)] d x \rightarrow \frac{1}{2}\|\psi\|^{2} \neq 0
$$

In fact, if we set $g_{n}=f_{n}^{2}=[\psi(x)]^{2} \sin ^{2}(n \pi x)$, then $g_{n} \rightarrow \frac{1}{2} \psi^{2}(x)$ because

$$
\begin{aligned}
\int g_{n}(x) \varphi(x) d x & =\int \psi^{2}(x) \sin ^{2}(n \pi x) \varphi(x) d x \\
& =\frac{1}{2} \int \psi^{2} \varphi d x-\frac{1}{2} \int \varphi^{2} \psi \cos (2 \pi n x) d x \\
& \rightarrow \frac{1}{2} \int \psi^{2} \varphi d x
\end{aligned}
$$

So $g_{n} \rightharpoonup \frac{1}{2} \psi^{2}$

Proof. (b)

$$
g_{n}= \begin{cases}n & 0<x<\frac{1}{n} \\ 0 & \text { otherwise }\end{cases}
$$

$\left\|g_{n}\right\|=\sqrt{n},\left(g_{n}\right)$ is unbounded, so $g_{n} \not \neg g$. In fact, $g_{n} \rightharpoonup \delta \in \mathcal{D}^{\prime}(\mathbb{R})$.

$$
h_{n}=\left\{\begin{array}{rl}
n^{1 / 4} & 0<x<\frac{1}{n} \\
0 & \text { otherwise }
\end{array}\right.
$$

$\left\|h_{n}\right\|=0$, and $\left(h_{n}\right)$ is strongly and weakly convergent to $0 . \frac{1}{2}$ is the critical value for $L^{2}$, and $\frac{1}{p}$ is the critical value for $L^{p}$.

Remark 8.52.

If $f_{n} \rightharpoonup f$ and $f_{n} \rightarrow g$, then we must have $f=g$ because

$$
\begin{aligned}
& \left\langle f_{n}, h\right\rangle \rightarrow\langle f, h\rangle \forall h \in \mathcal{H} \\
& \left\langle f_{n}, h\right\rangle \rightarrow\langle g, h\rangle \forall h \in \mathcal{H}
\end{aligned}
$$

Since $\langle f, h\rangle=\langle g, h\rangle \forall h$, we have that $f=g$.

### 8.7 The Banach-Alaoglu Theorem

Definition 8.53. Weakly Sequentially Compact
page 208 and Notes 2/23/11

A set $K \subset \mathcal{H}$ is weakly sequentially compact if for any sequence $\left(x_{n}\right) \subset K$ there exists a subsequence $\left(x_{n_{k}}\right)$ such that $x_{n_{k}} \rightharpoonup x \in K$.

Theorem 8.54. Banach-Alaoglu Theorem
page 208 and Notes $2 / 23 / 11$

Suppose that $\mathcal{H}$ is a separable Hilbert space and $\bar{B}=\{x \in \mathcal{H} \mid\|x\| \leq 1\}$ is the closed unit ball. Then $\bar{B}$ is weakly sequentially compact.

## Remarks

1. $\bar{B}$ is not strongly compact if $\mathcal{H}$ is infinite-dimensional. Ex: $\left\{e_{n}\right\}$ is an orthonormal basis, but $\left(e_{n}\right)$ has no convergent subsequence
2. This can be thought of as a replacement of the Heine-Borel theorem in the infinite-dimensional case

Proof. Let $\left\{y_{k} \mid k \in \mathbb{N}\right\}$ be a dense subset of $\mathcal{H}$. Consider $\left(\left\langle x_{n}, y_{1}\right\rangle\right)_{n} \subset \mathbb{C}$. By Cauchy-Schwarz, $\mid\left\langle x_{n}, y\right\rangle \leq\left\|x_{n}\right\|\left\|y_{1}\right\| \leq\left\|y_{1}\right\|$, so the sequence is bounded, and thus there exists a subsequence of $\left(x_{n}\right)$, denoted $\left(x_{n, 1, k}\right)_{k}=\left(x_{1, k}\right)$ such that $\left\langle x_{1, k}, y_{1}\right\rangle$ converges as $k \rightarrow \infty$. Pick a subsequence $\left(x_{2, k}\right)$ of $\left(x_{1, k}\right)$ such that $\left\langle x_{2, k}, y_{2}\right\rangle$ converges as $k \rightarrow \infty$. Let $x_{j}=x_{j, j}$ be the diagonal sequence. Then $\left\langle x_{j}, y_{n}\right\rangle$ converges for every $y_{k}$ as $j \rightarrow \infty$ in this dense subset of $\mathcal{H}$. This defines a bounded linear functional $F$ on $D=\left\{y_{k} \mid k \in \mathbb{N}\right\}$. By the Bounded Linear Transformation Theorem, this extends to a bounded linear functional $\bar{F}: \mathcal{H} \rightarrow \mathbb{C}$ such that $\bar{F}\left(y_{k}\right)=\lim _{j \rightarrow \infty}\left\langle x_{j}, y_{k}\right\rangle$ for all $k \in \mathbb{N}$. By the Riesz Representation Theorem, there exists $x \in \mathcal{H}$ such that $\left\langle x, y_{k}\right\rangle=\lim _{j \rightarrow \infty}\left\langle x_{j}, y_{k}\right\rangle$ for all $k \in \mathbb{N}$. Since $\left\{y_{k}\right\}$ is dense in $\mathcal{H}$ and $\|x\| \leq 1,\langle x, y\rangle=\lim _{j \rightarrow \infty}\left\langle x_{j}, y\right\rangle$ for all $y \in \mathcal{H}$, and thus $x_{j} \rightharpoonup x .\|x\| \leq \liminf _{j \rightarrow \infty}\left\|x_{j}\right\| \leq 1$, so $x \in \bar{B}$.

## Remark 8.55.

Notes 2/23/11

1. We don't need $\mathcal{H}$ to be separable (restrict to a closed subspace spanned by $\left\{x_{n}\right\}$ which is separable)
2. Generalization to Banach spaces: the unit ball of $X^{*}$ is weak-* compact (equivalent to being weak compact if $X$ is reflexive, i.e. $X^{* *}=X$ )

## Definition 8.56. Weakly Sequentially Closed

Notes 2/23/11

A set $F \subset \mathcal{H}$ is weakly sequentially closed if whenever $\left(x_{n}\right) \subset F$ is a sequence and $x_{n} \rightharpoonup x$, then $x \in F$.

Weakly closed implies strongly closed, but not conversely if $\mathcal{H}$ is infinite-dimensional. For example, let

$$
\begin{aligned}
S & =\{x \in \mathcal{H} \mid\|x\|=1\} \\
\bar{B} & =\{x \in \mathcal{H} \mid\|x\| \leq 1\}
\end{aligned}
$$

$S$ is not weakly closed because $\left(e_{n}\right) \subset S, e_{n} \rightharpoonup 0 \notin S$. $\bar{B}$ is weakly closed because if $x_{n} \rightharpoonup x$, then $\|x\| \leq \liminf \left\|x_{n}\right\|$. The weak closure of $S$ is $\bar{B}$.

## Definition 8.58. Weakly Sequentially Lower Semicontinuous

page 208 and Notes $2 / 23 / 11$

A function $f: D \subset \mathcal{H} \rightarrow \mathbb{R}$ is weakly sequentially lower semicontinuous if

$$
x_{n} \rightharpoonup x \quad \Rightarrow \quad f(x) \leq \liminf _{n \rightarrow \infty} f\left(x_{n}\right)
$$

Example: $\|\cdot\|: \mathcal{H} \rightarrow \mathbb{R}$ is weakly sequentially lower semicontinuous.

## Remark 8.59.

Notes 2/23/11

Weakly sequentially lower semicontinuous implies strongly sequentially lower semicontinuous, but not conversely.

## Theorem 8.60.

page 209 and Notes $2 / 23 / 11$

Suppose that $D$ is a weakly closed, bounded (in norm) subset in a Hilbert space $\mathcal{H}$ and $f: D \rightarrow \mathbb{R}$ is a weakly sequentially lower semicontinuous function. Then $f$ is bounded from below $\left(m=\inf _{x \in D} f(x)>-\infty\right)$ and there exists $x \in D$ such that $f(x)=m$.

### 8.8 Chapter Summary

We begin by defining what it means for a bounded linear operator $P$ to be a projection (with "opposite" $Q=I-P$ ), and we explore relationship between projections and direct sum decompositions: $P$ a projection $\Leftrightarrow X=\operatorname{ran} P \oplus \operatorname{ker} P$. We introduce orthogonal projections and show that they are bounded and self-adjoint. We explore the connection between orthogonal projections $P(\Rightarrow \mathcal{H}=\operatorname{ran} P \oplus \operatorname{ker} P)$ and direct sum decompositions ( $\mathcal{M}$ closed) $\mathcal{H}=\mathcal{M} \oplus \mathcal{M}^{\perp}\left(\Rightarrow P\right.$, ran $P=\mathcal{M}$, ker $\left.P=\mathcal{M}^{\perp}\right)$.

Recall from Chapter 5 that a linear functional is bounded iff it is continuous. We introduce the Riesz Representation Theorem: for all $\varphi \in \mathcal{H}^{*}$, there exists $y \in \mathcal{H}$ such that $\varphi(x)=\langle y, x\rangle$. This gives us that all Hilbert spaces are self-dual: $\mathcal{H}^{* *}=\mathcal{H}$. This is because the map $J_{1}: \mathcal{H} \rightarrow \mathcal{H}^{*}$ defined by $J_{1} y=\varphi_{y}$ identifies $\mathcal{H}$ with its dual space, $\mathcal{H}^{*}$. Similarly, we can define a map $J_{2}$ that identifies $\mathcal{H}^{*}$ with its dual space, $\mathcal{H}^{* *}$.

Thus, $\mathcal{H}$ and $\mathcal{H}^{* *}\left(\right.$ and $\left.\mathcal{H}^{*}\right)$ have the same cardinality. And since we know (Chapter 5) that for every $x \in \mathcal{H}$ we can define a functional $F_{x} \in \mathcal{H}^{* *}$ by $F_{x}(\varphi)=\varphi(x)$, we therefore know that all linear functionals in $\mathcal{H}^{* *}$ are of this form.

We use the Riesz Representation Theorem to prove the existence of the adjoint of a bounded operator on a Hilbert space: $\langle x, A y\rangle=\left\langle A^{*} x, y\right\rangle$. Examples:

- Matrix: $A^{*}=A^{T}\left(\overline{A^{T}}\right.$ if $A$ is complex)

$$
-\langle x, A y\rangle=x^{T} A y,\left\langle A^{*} x, y\right\rangle=\left(A^{*} x\right)^{T} y=x^{T}\left(A^{*}\right)^{T} y
$$

- Integral operator $K f(x)=\int_{0}^{1} k(x, y) f(y) d y: K^{*} f(x)=\int_{0}^{1} \overline{k(y, x)} f(y) d y$
- Shift operators: $S^{*}=T, T^{*}=S$

We verify that for a bounded linear operator $A$, a solvability condition for $A x=y$ is that $\langle y, z\rangle=0$ for all $z \in \operatorname{ker} A^{*} \Leftrightarrow \operatorname{ran} A \subset\left(\operatorname{ker} A^{*}\right)^{\perp}$. We use this fact to prove that for a bounded linear operator $A$,

$$
\overline{\operatorname{ran} A}=\left(\operatorname{ker} A^{*}\right)^{\perp}, \quad \operatorname{ker} A=\left(\operatorname{ran} A^{*}\right)^{\perp} .
$$

Equivalently,

$$
\mathcal{H}=\underbrace{\left(\operatorname{ker} A^{*}\right)^{\perp}}_{\operatorname{ran} A} \oplus \underbrace{(\operatorname{ran} A)^{\perp}}_{\text {ker } A^{*}} .
$$

Next we have some definitions. We define what it means for a bounded linear operator to be self-adjoint, and we prove that for a bounded self-adjoint operator $A$,

$$
\|A\|=\sup _{\|x\|=1}|\langle x, A x\rangle|, \quad\left\|A^{*} A\right\|=\|A\|^{2}
$$

Examples:

- A matrix is self-adjoint if it is symmetric (or Hermitian, if it is complex).
- An integral operator $K f(x)=\int_{0}^{1} k(x, y) f(y) d y$ is self-adjoint if $k(x, y)=\overline{k(y, x)}$

We say that an operator is unitary/orthogonal if it is invertible and $\langle U x, U y\rangle_{\mathcal{H}_{2}}=\langle x, y\rangle_{\mathcal{H}_{1}} \Leftrightarrow U^{*} U=$ $U U^{*}=I$. We say that an operator is normal if $T^{*} T=T T^{*}$. (Self-adjoing and unitary operators are normal.)

Now we revisit weak convergence. For Hilbert spaces, the Riesz Representation Theorem gives us an equivalent definition: $x_{n} \rightharpoonup x$ if $\left\langle x_{n}, y\right\rangle \rightarrow\langle x, y\rangle \forall y \in \mathcal{H} \Leftrightarrow \varphi\left(x_{n}\right) \rightarrow \varphi(x) \forall \varphi \in \mathcal{H}^{*}$. We mention 3 reasons why a sequence may converge weakly but not strongly: oscillation, concentration, and escape to infinity. We prove that for a weakly convergent sequence $\left(x_{n}\right),\|x\| \leq \lim \inf \left\|x_{n}\right\|$. We also prove that if $\lim \left\|x_{n}\right\|=\|x\|$, then $\left(x_{n}\right)$ converges to $x$ strongly. The Banach-Alaoglu Theorem tells us that the closed unit ball of a Hilbert space is weakly compact.

We define what it means for a function to be convex, and we say a few words about lower semicontinuous functions. We finish the chapter with Mazur's Theorem, which tells us that if $x_{n} \rightharpoonup x$, then there exists a sequence $\left(y_{n}\right)$ of finite convex combinations of $\left\{x_{n}\right\}$ that converges strongly to $x$.

## 9 The Spectrum of Bounded Linear Operators

### 9.0 Introduction

Remark 9.1.
page 215 and Notes $3 / 2 / 11$

Consider the following initial boundary value problem for a variable coefficient, linear equation:

$$
\begin{array}{lr}
u_{t}=u_{x x}-q(x) u & 0<x<1, t>0 \\
u(0, t)=0, u(1, t)=0 & t \geq 0 \\
u(x, 0)=f(x) & 0 \leq x \leq 1
\end{array}
$$

Using separation of variables, we assume

$$
u(x, t)=\sum_{n=1}^{\infty} a_{n}(t) u_{n}(x)
$$

where $\left\{u_{n} \mid n \in \mathbb{N}\right\}$ is an orthonormal basis of $L^{2}([0,1])$. We find that

$$
\frac{d a_{n}}{d t}=-\lambda_{n} a_{n} n
$$

and the $u_{n}$ satisfy

$$
-\frac{d^{2} u_{n}}{d x^{2}}+q u_{n}=\lambda_{n} u_{n}
$$

Then the $u_{n}$ are eigenvectors of the linear operator $A$. Thus, $A u_{n}=\lambda_{n} u_{n}$, where $A$ is defined by

$$
A u=-\frac{d^{2} u}{d x^{2}}+q u
$$

We want a complete set of eigenvectors of $A$, or equivalently, to diagonalize $A$. This is an example of what we do in spectral theory.

### 9.1 Diagonalization of Matrices

## Remark 9.2.

page 218 and Notes $3 / 2 / 11$

The concept of the spectrum of an operator on a Banach/Hilbert space is a generalization of eigenvalues for matrices. Let $A \in \mathcal{B}(X)$. When $\operatorname{dim} X<\infty$ then we can identify it with a a matrix $\tilde{A}$. For any $\lambda \in \mathbb{C}$ we have two possibilities:

1. $\lambda I-A$ is nonsingular $\Leftrightarrow \operatorname{det}(\lambda I-A)=0 \Leftrightarrow(\lambda I-A)^{-1}$ exists
2. $\lambda I-A$ is singular $\Leftrightarrow$ there exists $x_{0}$ such that $(\lambda I-A) x_{0}=0$. Thus, $A x_{0}=\lambda x_{0}, \lambda$ is an eigenvalue, and $x_{0}$ is an eigenvector.

What happens if $\operatorname{dim} X=\infty ? ? ?$

### 9.2 The Spectrum

Definition 9.3. Resolvent Set
page 218 and Notes $3 / 2 / 11$

The resolvent set of a bounded operator $A$ on a Banach space $X$ is the set

$$
\begin{aligned}
\qquad(A) & =\{\lambda \in \mathbb{C} \mid(\lambda I-A) \text { is invertible }\} \\
\text { (by the bounded inverse theorem) } & =\{\lambda \in \mathbb{C} \mid(\lambda I-A) \in \mathcal{B}(X)\} \\
& =\{\lambda \in \mathbb{C} \mid(\lambda I-A) \text { is 1-1 and onto }\}
\end{aligned}
$$

Definition 9.4. Spectrum
page 218 and Notes $3 / 2 / 11$

The spectrum of $A$ is the set

$$
\begin{aligned}
\sigma(A) & =\mathbb{C} \backslash \rho(A) \\
& =\{\lambda \in \mathbb{C} \mid(\lambda I-A) \text { is not invertible }\}
\end{aligned}
$$

## Definition 9.5. Point Spectrum, Continuous Spectrum, Residual Spectrum

 page 219 and Notes $3 / 2 / 11$In general, $\sigma(A)$ can be expressed as $\sigma(A)=\sigma_{p}(A) \cup \sigma_{c}(A) \cup \sigma_{r}(A)$, where

1. $\sigma_{p}(A)=\{\lambda \in \mathbb{C} \mid(\lambda I-A)$ is not 1-1 $\}$
$\sigma_{p}(A)$ is called the point spectrum of $A$. In this case, since $(\lambda I-A)$ is not 1-1, there exists $x_{0} \in \operatorname{ker}(\lambda I-A)$ such that $(\lambda I-A) x_{0}=0 \Leftrightarrow A x_{0}=\lambda x_{0}$
2. $\sigma_{c}(A)=\{\lambda \in \mathbb{C} \mid(\lambda I-A)$ is 1-1 but not onto and $\overline{\operatorname{ran}(\lambda I-A)}=X\}$
$\sigma_{c}(A)$ is called the continuous spectrum of $A$
3. $\sigma_{r}(A)=\{\lambda \in \mathbb{C} \mid(\lambda I-A)$ is 1-1 but not onto and $\overline{\operatorname{ran}(\lambda I-A)} \neq X\}$
$\sigma_{r}(A)$ is called the residual spectrum of $A$

Example 9.6. Point, Continuous, and Residual Spectra Examples
Notes 3/7/11

1. A matrix on $\mathbb{C}^{n}$ has pure point spectrum
2. $M: L^{2}([0,1]) \rightarrow L^{2}([0,1]), f \mapsto x f, \sigma(M)=[0,1]$ has pure continuous spectrum
3. Consider the right shift operator $S$ on $\ell^{2}(\mathbb{N})$. $\lambda=0$ is in the residual spectrum

## Example 9.7.

Notes 3/2/11

Consider the Banach space $X=C([0,1])$ with the $\|\cdot\|_{\infty}$ norm. Define $A: X \rightarrow X$ by $A f(x)=x f(x)$. The boundedness of $A$ follows exactly as in HW7 (even though $X=L^{2}([0,1])$ on the HW, since we can take $\sup x=1)$. Find $\sigma(A)$. Claim: $\sigma(A)=\sigma_{r}(A)=[0,1]$.

For any $\lambda \in \mathbb{C}, f \in C([0,1])$, we have

$$
(\lambda I-A) f(x)=(\lambda-x) f(x)=0
$$

If $\lambda \neq x$ then $f(x)=0$. If $\lambda \notin[0,1]$ then $\sigma_{p}=\emptyset$.
For all $\lambda \notin[0,1]$, is $(\lambda I-A)$ onto? For every $g \in C([0,1])$, we want $f$ such that $f(x)(\lambda-x)=$ $g(x) \Rightarrow f(x)=\frac{g(x)}{\lambda-x} \in C([0,1])$, since $\lambda \notin[0,1]$ implies that $\lambda-x \neq 0 \forall x \in[0,1]$. Thus, $(\lambda I-A)$ is onto, and we can conclude that $\sigma(A) \subseteq[0,1]$.
It will be enough to prove the claim to show that $[0,1] \subseteq \sigma_{r}(A)$. Why? $[0,1] \subseteq \sigma_{r}(A) \subseteq \sigma(A) \subseteq[0,1]$.
Pick $\lambda \in[0,1]$. For every $g \in \operatorname{ran}(\lambda I-A)$ we have that

$$
\begin{aligned}
& g(x)=(\lambda-x) f(x) \text { for some } f \in X=C([0,1]) \\
& g(\lambda)=0
\end{aligned}
$$

So $h(x)=1 \notin \operatorname{ran}(\lambda I-A)$, since $g(\lambda)=0 \neq 1$. Therefore $(\lambda I-A)$ is not onto.

If $h \in \overline{\operatorname{ran}(\lambda I-A)}$ then there exists $\left(g_{n}\right) \subset \operatorname{ran}(\lambda I-A)$ such that $g_{n} \rightarrow h . h\left(\lambda=\lim _{n \rightarrow \infty} g_{n}(\lambda)(\lambda I-\right.$ $A)=0$. Thus, $\mathbf{1} \notin \overline{\operatorname{ran}(\lambda I-A)}$, so $\lambda \in \sigma_{r}(A)$.

Example 9.8.
page 219

Example 9.5 on page 219

## Definition 9.9. Resolvent

page 220 and Notes $3 / 4 / 11$

For $\lambda \in \rho(A)$, we define the resolvent of $A$ at $\lambda$ to be

$$
R_{\lambda}=(\lambda I-A)^{-1}, \quad R_{\lambda}: \rho(A) \subset \mathbb{C} \rightarrow \mathcal{B}(\mathcal{H})
$$

Example 9.10. Neumann Series
page 220 and Notes $3 / 4 / 11$

If $\|A\|<1$ then $(I-A)$ is invertible and

$$
(I-A)^{-1}=I+A+A^{2}+\ldots
$$

To show this, we define the partial sum:

$$
S_{N}=I+A+A^{2}+\ldots+A^{N}
$$

Next, we show that the sequence of partial sums is Cauchy:

$$
\begin{aligned}
\left\|A^{M+1}+\ldots+A^{N}\right\| & \leq\left\|A^{M+1}\right\|+\ldots+\left\|A^{N}\right\| \leq\|A\|^{M+1}+\ldots+\|A\|^{N} \\
& \leq \sum_{n=M+1}^{N}\|A\|^{n}
\end{aligned}
$$

$\sum_{n=1}^{\infty}<\infty$ if $\|A\|<1$, so the partial sums are Cauchy. Thus, $\sum_{n=0}^{\infty} A^{n}$ is Cauchy in $\mathcal{B}(\mathcal{H})$, and it converges since $\mathcal{B}(\mathcal{H})$ is complete.
(See Remark 9.12.)

## Example 9.11.

Notes 3/4/11

1. If $|\lambda|>\|A\|$ then $\lambda \in \rho(A)$
$(\lambda I-A)^{-1}=\left[\lambda\left(I-\frac{A}{\lambda}\right)\right]^{-1}=\frac{1}{\lambda}\left(I-\frac{A}{\lambda}\right)^{-1}$
$\uparrow$ this exists if $\|A / \lambda\|<1 \Rightarrow\|A\|<|\lambda|$
2. The resolvent set $\rho(A)$ is open in $\mathbb{C}$

Suppose $\lambda_{0} \in \rho(A)$. We write:

$$
\begin{aligned}
& (\lambda I-A)=\lambda_{0} I-A+\left(\lambda-\lambda_{0}\right) I=\left(\lambda_{0} I-A\right)\left[I+\left(\lambda-\lambda_{0}\right)\left(\lambda_{0} I-A\right)^{-1}\right] \\
& (\lambda I-A)^{-1}=\left[I+\left(\lambda-\lambda_{0}\right)\left(\lambda_{0} I-A\right)^{-1}\right]^{-1}\left(\lambda_{0} I-A\right)^{-1} \\
& \quad \uparrow \text { exists if }\left|\lambda-\lambda_{0}\right|<\frac{1}{\left\|\left(\lambda_{0} I-A\right)^{-1}\right\|}
\end{aligned}
$$

3. $R_{\lambda}: \lambda \mapsto(\lambda I-A)^{-1}$
$R_{\lambda}$ is an operator-valued analytic function on the open set $\rho(A) \subset \mathbb{C}$
4. $\sigma(A) \neq \emptyset$

## Remark 9.12.

Notes 3/4/11

In Example 9.10, it is not necessary that $\|A\|<1$ for $(I-A)^{-1}=I+A+A^{2}+\ldots$ to converge.
Rather, we require that $\lim _{n \rightarrow \infty}\left\|A^{n}\right\|^{1 / n}<1$.
$r(A)=\sup \{|\lambda| \mid \lambda \in \sigma(A)\}$ is the spectral redius of $A$. This is the radius of the smallest disc in $\mathbb{C}$ centered at 0 that contains $\sigma(A)$. Also, $r(A) \leq\|A\|$.

Theorem 9.14.
page 220 and Notes $3 / 4 / 11$

$$
r(A)=\lim _{n \rightarrow \infty}\left\|A^{n}\right\|^{1 / n} \quad \text { (and the limit exists) }
$$

Proof. Let $a_{n}=\log \left\|A^{n}\right\|$. (If $\left\|A^{n}\right\|=0$ for some $n$, i.e. $A$ is nilpotent, then $r(A)=0$.) Then

$$
\begin{aligned}
a_{m+n} & =\log \left\|A^{m+n}\right\| \\
& \leq \log \left\|A^{n}\right\|+\log \left\|A^{n}\right\| \\
& \leq a_{m}+a_{n} \quad \text { (subadditive) }
\end{aligned}
$$

We want to show that $\lim _{n \rightarrow \infty} \frac{a_{n}}{n}$ exists, where $\frac{a_{n}}{n}=\log \left\|A^{n}\right\|^{1 / n}$. Fix $n, m$ and write $n=m p+q$ with $0 \leq q<m$. Then we have

$$
\begin{aligned}
& a_{n}=a_{m p+q} \leq a_{m p}+a_{q} \\
& \frac{a_{n}}{n} \leq \frac{a_{m p}}{n}+\frac{a_{q}}{n}
\end{aligned}
$$

Note that $a_{m p} \leq p a_{m}$. Let $n \rightarrow \infty$ with $m$ fixed. Then $\frac{p}{n} \rightarrow \frac{1}{m}$ as $n \rightarrow \infty$, and

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \frac{a_{n}}{n} \leq \frac{a_{m}}{m} \tag{9.1}
\end{equation*}
$$

Taking the limit of (9.1) as $m \rightarrow \infty$, we obtain

$$
\limsup _{n \rightarrow \infty} \frac{a_{n}}{n} \leq \liminf _{m \rightarrow \infty} \frac{a_{m}}{m}
$$

So $\limsup _{n \rightarrow \infty} \frac{a_{n}}{n}=\limsup _{n \rightarrow \infty} \frac{a_{n}}{n}$, and the sequence converges.

## Example 9.15. Example for Theorem 9.14

Notes 3/4/11

$$
\begin{aligned}
A & =\mu I \\
\lambda I-A & =(\lambda-\mu) I \\
\sigma(A) & =\mu
\end{aligned}
$$

$$
\begin{aligned}
\|A\| & =|\mu|=r(A) \\
\left\|A^{n}\right\|^{1 / n} & =|\mu|
\end{aligned}
$$

Corollary 9.16.
page 221 and Notes $3 / 4 / 11$

If $A$ is self-adjoint then $r(A)=\|A\|$.

Proof. $\left\|A^{2}\right\|=\|A\|^{2}$ and $\left\|A^{2^{n}}\right\|=\|A\|^{2^{n}}$, so $\liminf _{n \rightarrow \infty}\left\|A^{n}\right\|^{1 / n}=\|A\|$ by taking the subsequence $n=2^{m}$.

### 9.3 The Spectral Theorem for Compact, Self-Adjoint Operators

### 9.3.1 Bounded, Self-Adjoint Operators

## Theorem 9.17.

page 222 and Notes $3 / 7 / 11$

If $A$ is bounded and self-adjoint, then every eigenvalue of $A$ is real and eigenvectors with different eigenvalues are orthogonal.

Related to Theorem 9.21.

Proof. If $A x=\lambda x$, then

$$
\begin{aligned}
& \langle x, A x\rangle=\langle x, \lambda x\rangle=\lambda\|x\|^{2} \\
& \langle A x, x\rangle=\langle\lambda x, x\rangle=\bar{\lambda}\|x\|^{2}
\end{aligned}
$$

If $A$ is self-adjoint (and $x \neq 0$ ), then $\lambda=\bar{\lambda} \Rightarrow \lambda \in \mathbb{R}$.
Case: $A$ has pure point spectrum.
If $A x=\lambda x$ and $A y=\mu y, x, y \neq 0, \lambda \neq \mu$, then

$$
\left.\begin{array}{l}
\langle x, A y\rangle=\mu\langle x, y\rangle \\
\langle A x, y\rangle=\bar{\lambda}\langle x, y\rangle=\lambda\langle x, y\rangle
\end{array}\right\} A=A^{*}, \text { so } \mu\langle x, y\rangle=\lambda\langle x, y\rangle
$$

If $\lambda \neq \mu$, then $\langle x, y\rangle=0$, i.e. $x \perp y$.
What about the continuous and residual spectra?

$$
\begin{aligned}
\|(A-\lambda I) x\|^{2} & =\langle(A-a I) x-i b x,(A-a I) x-i b x\rangle \quad \text { where } \lambda=a+i b \\
& =\langle(A-a I) x,(A-a I) x\rangle+\underline{\langle-i b x,(A-a I) x\rangle}+\underline{\langle(A-a I) x,-i b x\rangle}+\langle-i b x,-i b x\rangle \\
& =\|(A-a I) x\|^{2}+b^{2}\|x\|^{2} \\
& \geq b^{2}\|x\|^{2}
\end{aligned}
$$

Continuous Spectrum: See Proposition 9.18 and Remark 9.19.
Residual Spectrum: See Proposition 9.20.

## Proposition 9.18.

page 223 and Notes $3 / 7 / 11$

$$
|\operatorname{Im} \lambda| \cdot\|x\| \leq\|(A-a I) x\|
$$

## Remark 9.19.

Notes 3/7/11

Proposition 9.18 says that if $(A-\lambda I) x=y$, then $|\operatorname{Im} \lambda| \cdot\|x\| \leq\|y\|$. This means that if $\lambda \in \mathbb{R}$, we can estimate the solution, $x$, in terms of the RHS, $y$.

Applying this to the proof of Theorem 9.17, we see that if $\lambda \in \mathbb{C} \backslash \mathbb{R}$, it follows that
(a) $(A-\lambda I)$ is 1-1 because if $(A-\lambda I) x=0$ then $|\operatorname{Im} \lambda|\|x\|=0 \Rightarrow x=0$.
(b) $(A-\lambda I)$ has closed range. If $y_{n}=(A-\lambda I) x_{n}, y_{n} \in \operatorname{ran}(A-\lambda I), y_{n} \rightarrow y$, then we can bound

$$
\underbrace{\left\|x_{m}-x_{n}\right\|}_{\therefore \text { Cauchy }} \leq C \underbrace{\left\|y_{m}-y_{n}\right\|}_{\text {Cauchy }}
$$

So $x_{n} \rightarrow x,(A-\lambda I) x=y$, and $y \in \operatorname{ran}(A-\lambda I)$. So if $\lambda \in \mathbb{C} \backslash \mathbb{R}$, then $(A-\lambda I)$ is $1-1$ with closed range, so there is no complex-valued continuous spectrum.

## Proposition 9.20.

page 224 and Notes $3 / 7 / 11$

If $A$ is bounded and self-adjoint, then the residual spectrum is empty.

Proof. If $\lambda$ is in the residual spectrum, then there exists $y \in \mathcal{H}$ such that $\langle(A-\lambda I) x, y\rangle=0 \forall x \in \mathcal{H}$, so $y \perp \overline{\operatorname{ran}(A-\lambda I)}, y \neq 0$. Since $A$ is self-adjoint, $\langle x,(A-\bar{\lambda} I) y\rangle=0 \forall x \in \mathcal{H}$. This implies that $(A-\bar{\lambda} I) y=0$, so $y$ is an eigenvector of $A$ with eigenvalue $\bar{\lambda}$. We have 2 cases:

1. $\lambda \in \mathbb{C} \backslash \mathbb{R} \Rightarrow$ impossible ( $A$ has real eigenvalues)
2. $\lambda \in \mathbb{R}$. Then $\lambda$ is in the point and residual spectra $\Rightarrow$ impossible.

## Theorem 9.21.

page 223 and Notes $3 / 7 / 11$

If $A$ is a bounded, self-adjoint operator on a Hilbert space $\mathcal{H}$, then $\sigma(A)$ is real and contained in the interval $[-\|A\|,\|A\|]$. The residual spectrum is empty.

Related to Theorem 9.17.

## Proposition 9.22.

page 223

If $A$ is a bounded operator on a Hilbert space (not necessarily self-adjoint!) and $\lambda \in \sigma_{r}(A)$, then $\bar{\lambda} \in \sigma_{p}\left(A^{*}\right)$. In other words, $\sigma_{r}(A) \subseteq \sigma_{p}\left(A^{*}\right)$.

Bounded, self-adjoint operators have

- Spectral radius $r(A)=\|A\|$ (See Corollary 9.16)
- Real eigenvalues (See Theorem 9.17)
- Orthogonal eigenvectors (See Theorem 9.17)
- Empty residual spectrum (See Proposition 9.20)


### 9.3.2 Compact Operators

## Definition 9.24. Compact Operator

Notes 3/9/11
$K: \mathcal{H} \rightarrow \mathcal{H}, D \in \mathcal{B}(\mathcal{H})$ is compact if it maps bounded sets to precompact sets.

## Remark 9.25. Precompact

Notes 3/9/11

Remember: a set is precompact if it is bounded and "almost" finite-dimensional.

Example 9.26. The Hilbert Cube
page 230 and Notes 3/9/11

Let $\mathcal{H}=\ell^{2}(\mathbb{N})$. The Hilbert cube

$$
C=\left\{\left(x_{1}, x_{2}, \ldots, x_{n}, \ldots\right)| | x_{n} \left\lvert\, \leq \frac{1}{n}\right.\right\}
$$

is closed and precompact. Hence, $C$ is a compact subset of $\mathcal{H}$.

Example 9.27. Diagonal Operators Are Compact
page 230

The diagonal operator : $\ell^{2}(\mathbb{N}) \rightarrow \ell^{2}(\mathbb{N})$ defined by

$$
A\left(x_{1}, x_{2}, \ldots, x_{n}, \ldots\right)=\left(\lambda_{1} x_{1}, \lambda_{2} x_{2}, \ldots, \lambda_{n} x_{n}, \ldots\right)
$$

is compact iff $\lambda_{n} \rightarrow 0$ as $n \rightarrow \infty$.

## Example 9.28. Compactness of Operators

Notes 3/9/11

1. Any operator with finite $\operatorname{rank}(\operatorname{rank} A=\operatorname{dim} \operatorname{ran} A)$ is compact
2. $I: \mathcal{H} \rightarrow \mathcal{H}$ is not compact if $\operatorname{dim} \mathcal{H}=\infty$
3. $L^{2}([0,1]), K f(x)=\int_{0}^{x} f(y) d y$ is a compact operator. If $\|f\|_{L^{2}} \leq M$, then

$$
\left|\int_{0}^{x} f(y) d y\right| \leq \int_{0}^{1}|f(y)| d y \leq\left(\int_{0}^{1}|f(y)|^{2} d y\right)^{1 / 2} \leq M
$$

Define $F(x)=\int_{0}^{x} f(y) d y$. Then

$$
\left|F\left(x_{2}\right)-F\left(x_{1}\right)\right|=\left|\int_{x_{1}}^{x_{2}} f(y) d y\right| \leq\left(\int_{x_{1}}^{x_{2}} 1 \cdot d y\right)^{1 / 2}\left(\int_{x_{1}}^{x_{2}}|f(y)|^{2} d y\right)^{1 / 2} \leq M\left|x_{2}-x_{1}\right|^{1 / 2}
$$

$\{K f \mid\|f\| \leq M\}$ is bounded and equicontinuous. Thus, $H^{2}([0,1])$ is compactly embedded in $L^{2}([0,1])$. It follows that $\left\{K f \mid\|f\|_{L^{2}} \leq M\right\}$ is precompact in $C([0,1])$ by Arzela-Ascoli, so it is precompact in $L^{2}([0,1])$.

If $f(x)=\sum_{n=1}^{\infty} b_{n} \sin (n \pi x)$, then $K f(x)=\sum_{n=1}^{\infty} \frac{b_{n}}{n \pi}-\sum_{n=1}^{\infty} \frac{b_{n}}{n \pi} \cos (n \pi x)$

### 9.3.3 Compact, Self-Adjoint Operators

## Remark 9.29.

page 223 and Notes $3 / 9 / 11$

Given: $A: \mathcal{H} \rightarrow \mathcal{H}, A$ is compact and self-adjoint, $\mathcal{H}$ is a separable Hilbert space We will prove:

1. $A$ has at least one eigenvalue
2. If $A$ leaves a subspace $M \subset \mathcal{H}$ invariant $(A: M \rightarrow M)$, then $A$ leaves $M^{\perp}$ invariant, and $\mathcal{H}=M \oplus M^{\perp}$

Idea: if we have $A \varphi_{n}=\lambda_{n} \varphi_{n}$, then we can get the largest eigenvalue by maximizing $A\left(\sum c_{n} \varphi_{n}\right)=$ $\sum \lambda_{n} c_{n} \varphi_{n}$.

## Theorem 9.30.

page 225 and Notes $3 / 9 / 11$

Suppose $A: \mathcal{H} \rightarrow \mathcal{H}$ is compact and self-adjoint. Then $A$ has an eigenvector with eigenvalue $\lambda$ with $\lambda=\|A\|$ and $/$ or $\lambda=-\|A\|$.

Proof. Recall: since $A$ is self-adjoint, $\|A\|=\sup _{\|x\|=1}|\langle x, A x\rangle|$. Choose a sequence $\left(x_{n}\right) \subset \mathcal{H}$ with $\left\|x_{n}\right\|=1$
and $\left\langle x_{n}, A x_{n}\right\rangle \rightarrow \lambda$ as $n \rightarrow \infty, \lambda= \pm\|A\|$. Then we have

$$
\begin{array}{rll}
\left\|(A-\lambda I) x_{n}\right\|^{2} & = & \left\langle(A-\lambda I) x_{n},(A-\lambda I) x_{n}\right\rangle \\
& =\left\langle A x_{n}, A x_{n}\right\rangle-2 \lambda\left\langle x_{n}, A x_{n}\right\rangle+\lambda^{2}\left\langle x_{n}, x_{n}\right\rangle \\
= & \underbrace{\left\|A x_{n}\right\|^{2}}-2 \lambda\left\langle x_{n}, A x_{n}\right\rangle+\lambda^{2} \\
& \leq\|A\|^{2}\left\|x_{n}\right\|^{2}=\lambda^{2} & \\
\leq & 2 \lambda^{2}-2 \lambda\left\langle x_{n}, A x_{n}\right\rangle & \rightarrow 0 \text { as } n \rightarrow \infty
\end{array}
$$

So $(A-\lambda I) x_{n} \rightarrow 0$ as $n \rightarrow \infty$, and thus $x_{n}-\frac{1}{\lambda} A x_{n} \rightarrow 0$ (assuming $\lambda \neq 0$, in which case $\|A\|=0$ and everything is an eigenvalue). Since $\left(x_{n}\right)$ is bounded $\left(\left\|x_{n}\right\|=1 \forall n\right), A x_{n} \rightarrow y$ by the compactness of $A$. So $x_{n} \rightarrow \frac{y}{\lambda}$ and $(A-\lambda I) y=0 .\|y\|=\lambda \neq 0$, since $\left\|x_{n}\right\|=1$ and $x_{n} \rightarrow y$. So $A$ has eigenvector $y$ with eigenvalue $\lambda$.

## Proposition 9.31.

page 224 and Notes $3 / 9 / 11$

1. Any nonzero eigenvalue of a compact, self-adjoint operator has a finite multiplicity (multiplicity $\equiv$ the dimension of the eigenspace).
2. If $\lambda_{n}$ is a sequence of eigenvalues and $\lambda_{n} \rightarrow L$, then we must have that $L=0$.

## Theorem 9.32. Spectral Theorem for Compact, Self-Adjoint Operators

 page 225 and Notes $3 / 11 / 11$If $A: \mathcal{H} \rightarrow \mathcal{H}$ is a compact, self-adjoint operator on a Hilbert space $\mathcal{H}$ then there is a finite or countably infinite sequence $\left(\lambda_{n}\right)$ of nonzero real eigenvalues and orthogonal eigenvectors $\left(\varphi_{n}\right)$ such that

$$
A \varphi_{n}=\lambda_{n} \varphi_{n}
$$

where $\left|\lambda_{1}\right| \geq\left|\lambda_{2}\right| \geq \ldots . \lambda_{n} \rightarrow 0$ as $n \rightarrow \infty$ if there are infinitely many $\lambda_{n}$ 's and

$$
\begin{aligned}
A x & =\sum_{n} \lambda_{n}\left\langle\varphi_{n}, x\right\rangle \varphi_{n} \\
x & =\sum\left\langle\varphi_{n}, x\right\rangle \varphi_{n}+n \quad \text { where } n \in \operatorname{ker} A, \quad \text { ker } A \perp \underbrace{\left\langle\varphi_{n}\right\rangle}_{\text {span }}
\end{aligned}
$$

Let $P_{n}: \mathcal{H} \rightarrow \mathcal{H}$ be the orthogonal projection onto the eigenspace with eigenvalue $\lambda_{n}$ (eigenvectors of bounded, self-adjoint operators are orthogonal; see Theorem 9.17). Then

$$
A=\sum \lambda_{n} P_{n}
$$

We are representing $A$ as a sum of linear projections because $\lambda_{n} \rightarrow 0$, and so the sum converges uniformly.

Proof. To see that the sum converges uniformly to $A$, we compute

$$
\left\|A x-\sum_{n=1}^{N} \lambda_{n} P_{n} x\right\|=\sum_{n=N+1}^{\infty}\left|\lambda_{n}\left\langle\varphi_{n}, x\right\rangle \varphi_{n}\right|^{2} \leq\left|\lambda_{N+1}\right|^{2}\|x\|^{2}
$$

Also, if we let $P_{0}$ be the orthogonal projection onto $\operatorname{ker} A$, then

$$
P_{0}+\sum P_{n}=I
$$

is strongly convergent. This is an example of what's called "resolution of the identity." Note that the $\lambda_{i}$ 's gave us uniform convergence above. For bounded (and unbounded) self-adjoint operators with continuous spectrum we need to use resolutions of identity that involve integrals (instead of sums).

### 9.4 Functions of Operators $=$ Functional Calculus

## Definition 9.33. Function of an Operator

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If $f: \sigma(A) \subset \mathbb{C} \rightarrow \mathbb{C}$ is a bounded function, then we define

$$
f(A)=\sum f\left(\lambda_{n}\right) P_{n}+f(0) P_{0}
$$

- $f$ is uniformly convergent if $f\left(\lambda_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$
- $f$ is strongly convergent if $f\left(\lambda_{n}\right) \nrightarrow 0$ as $n \rightarrow \infty$

Note that $\sigma(A)=\left\{\lambda_{n}\right\} \cup\{0\}$ if $\operatorname{dim} H=\infty$

- If there are finitely many $\lambda_{n}$, then $0 \in \sigma_{p}(A)$
- If there are countably many $\lambda_{n}$, then $0 \in \sigma_{c}(A)$


## Example 9.34

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Suppose $A$ is a positive (see Definition 8.32), self-adjoint compact operator. Then

$$
\langle x, A x\rangle \geq 0 \quad \text { implies } \quad \lambda_{n} \geq 0 \forall n
$$

We can define the positive square root of $A$ as

$$
\begin{aligned}
\sqrt{A} & =\sum \lambda_{n}^{1 / 2} P_{n} \\
(\sqrt{A})^{2} & =\sum \lambda_{n} P_{n}=A
\end{aligned}
$$

In general, if $A$ is compact then
$T=A^{*} A$ is positive and self-adjoint because $\quad\langle x, T x\rangle=\left\langle x, A^{*} A x\right\rangle=\langle A x, A x\rangle \geq 0$

$$
\sqrt{T}=|A|, \quad|A|^{2}=T=A^{*} A
$$

## Definition 9.35. Polar Decomposition

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$A=U|A|$, where $U: \operatorname{ran}|A| \rightarrow \operatorname{Im} A$ is a unitary operator

## Definition 9.36. Fredholm Operator, Index

Notes 3/11/11

A bounded operator $A: \mathcal{H} \rightarrow \mathcal{H}$ is Fredholm if
(a) $\operatorname{ran} A$ is closed
(b) $\operatorname{dim} \operatorname{ker} A$ is finite
(c) $\operatorname{codim} \operatorname{ran} A$ is finite $\Leftrightarrow \operatorname{dim} \operatorname{ker} A^{*}$ is finite

- codim $\operatorname{ran} A=\operatorname{dim} \operatorname{ker} A^{*}$ (recall that $\mathcal{H}=\operatorname{ran} A \oplus \operatorname{ker} A^{*}$ when $\operatorname{ran} A$ is closed)

We define the index by

$$
\text { index } A=\operatorname{dim}(\operatorname{ker} A)-\operatorname{codim}(\operatorname{ran} A)=\operatorname{dim}(\operatorname{ker} A)-\operatorname{dim}\left(\operatorname{ker} A^{*}\right)
$$

## Example 9.37. Fredholm or not?

Notes 3/11/11
(a) $I$ is Fredholm with index $=0$
(b) $A\left(x_{1}, x_{2}, x_{3}, \ldots\right) \mapsto\left(x_{1}, \frac{1}{2} x_{2}, \frac{1}{3} x_{3}, \ldots\right)$ is not Fredholm because the range is not closed
(c) The right shift operator, $S$, is Fredholm with index $=-1$

If $A$ is Fredholm with $\operatorname{index}(A)=0$ then we have Fredholm alternative for solving the equation $A x=y$, and there are 2 possibilities:

1. $A$ is one-to-one and we can solve the equation for every $y \in \mathcal{H}$
2. $A$ is not one-to-one, and we can only solve the equation if $y \perp \operatorname{ker} A^{*}$

## Theorem 9.38. Riesz-Schauder Theorem

Notes 3/11/11

If $K$ is a compact, self-adjoint operator and $\lambda \neq 0$ then $A=\lambda I-K$ is Fredholm with index 0 .

### 9.5 Chapter Summary

$$
\begin{array}{rll}
U^{*} A U & =U^{*}(A U)=U^{*}\left(A\left[\begin{array}{lll}
u_{1} & u_{2} & \cdots \\
u_{k}
\end{array}\right]\right) \\
& =U^{*}\left[\begin{array}{llll}
A u_{1} & A u_{2} & \cdots & A u_{k}
\end{array}\right] \\
& =U^{*}\left[\begin{array}{llll}
\lambda_{1} u_{1} & \lambda_{2} u_{2} & \cdots & \lambda_{k} u_{k}
\end{array}\right] \\
& =\left[\begin{array}{llll}
\lambda_{1} e_{1} & \lambda_{2} e_{2} & \cdots & \lambda_{k} e_{k}
\end{array}\right] & \text { (because } \left.U^{*} u_{k}=e_{k}\right) \\
& =\left[\begin{array}{cccc}
\lambda_{1} & 0 & \cdots & 0 \\
0 & \lambda_{2} & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & \lambda_{k}
\end{array}\right]=D
\end{array}
$$

| Operator | Spectrum | Point | Continuous | Residual |
| :---: | :---: | :---: | :---: | :---: |
| Bounded, Linear | Closed \& Nonempty, $r(A)=\lim \left\\|A^{n}\right\\|^{1 / n}$ |  |  | $\begin{gathered} \lambda \in \sigma_{r}(A) \Rightarrow \\ \bar{\lambda} \in \sigma_{p}\left(A^{*}\right) \end{gathered}$ |
| Bounded, Self-Adjoint | $\begin{gathered} \sigma(A) \subset[-\\|A\\|,\\|A\\|] \\ r(A)=\\|A\\| \\ \hline \end{gathered}$ | real | real | empty |
| Compact, Self-Adjoint |  | $\begin{gathered} -\\|A\\| \in \sigma_{p}(A) \text { or } \\ \\|A\\| \in \sigma_{p}(A) \end{gathered}$ | $\begin{gathered} \sigma_{c}(A)=\{0\} \text { or } \\ \sigma_{c}(A)=\emptyset \\ \hline \end{gathered}$ | empty |

